## Quantum Field Theory I + II



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## Literature

This is a writeup of my Master programme course on Quantum Field Theory I (Chapters 1-6) and Quantum Field Theory II. The primary source for this course has been

- Peskin, Schröder: An introduction to Quantum Field Theory, ABP 1995,
- Itzykson, Zuber: Quantum Field Theory, Dover 1980,
- Kugo: Eichtheorie, Springer 1997,
which I urgently recommend for more details and for the many topics which time constraints have forced me to abbreviate or even to omit. Among the many other excellent textbooks on Quantum Field Theory I particularly recommend
- Weinberg: Quantum Field Theory I + II, Cambridge 1995,
- Srednicki: Quantum Field Theory, Cambridge 2007,
- Banks: Modern Quantum Field Theory, Cambridge 2008
as further reading. All three of them oftentimes take an approach different to the one of this course. Excellent lecture notes available online include
- A. Hebecker: Quantum Field Theory,
- D. Tong: Quantum Field Theory.

Special thanks to Robert Reischke ${ }^{1}$ for his fantastic work in typing these notes.

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## Chapter 1

## The free scalar field

### 1.1 Why Quantum Field Theory?

In (non-relativistic) Quantum Mechanics, the dynamics of a particle is described by the time-evolution of its associated wave-function $\psi(t, \vec{x})$ with respect to the non-relativistic Schrödinger equation

$$
\begin{equation*}
i \hbar \frac{\partial}{\partial t} \psi(t, \vec{x})=H \psi(t, \vec{x}), \tag{1.1}
\end{equation*}
$$

with the Hamilitonian given by $H=\frac{\hat{\bar{p}}^{2}}{2 m}+V(\hat{x})$. In order to achieve a Lorentz invariant framework, a naive approach would start by replacing this non-relativistic form of the Hamiltonian by a relativistic expression such as

$$
\begin{equation*}
H=\sqrt{c^{2} \hat{\vec{p}}^{2}+m^{2} c^{4}} \tag{1.2}
\end{equation*}
$$

or, even better, by modifying the Schrödinger equation altogether such as to make it symmetric in $\frac{\partial}{\partial t}$ and the spatial derivative $\vec{\nabla}$. However, the central insight underlying the formulation of Quantum Field Theory is that this is not sufficient. Rather, combining the principles of Lorentz invariance and Quantum Theory requires abandoning the single-particle approach of Quantum Mechanics.

- In any relativistic Quantum Theory, particle number need not be conserved, since the relativistic dispersion relation $E^{2}=c^{2} \vec{p}^{2}+m^{2} c^{4}$ implies that energy can be converted into particles and vice versa. This requires a multi-particle framework.
- Unitarity and causality cannot be combined in a single-particle approach: In Quantum Mechanics, the probability amplitude for a particle to propagate from position $\vec{x}$ to $\vec{y}$ is

$$
\begin{equation*}
G(\vec{x}, \vec{y})=\langle\vec{y}| e^{-\frac{i}{\hbar} H t}|\vec{x}\rangle . \tag{1.3}
\end{equation*}
$$

One can show that e.g. for the free non-relativistic Hamiltonian $H=\frac{\hat{p}^{2}}{2 m}$ this is non-zero even if $x^{\mu}=\left(x^{0}, \vec{x}\right)$ and $y^{\mu}=\left(y^{0}, \vec{y}\right)$ are at a spacelike distance. The problem persists if we replace $H$ by a relativistic expression such as (1.2).

Quantum Field Theory (QFT) solves both these problems by a radical change of perspective:

- The fundamental entities are not the particles, but the field, an abstract object that penetrates spacetime.
- Particles are the excitations of the field.

Before developing the notion of an abstract field let us try to gain some intuition in terms of a mechanical model of a field. To this end we consider a mechanical string of length $L$ and tension $T$ along the $x$-axis and excite this string in the transverse direction. Let $\phi(x, t)$ denote the transverse excitation of the string. In this simple picture $\phi(x, t)$ is our model for the field. This system arises as the continuum limit of $N$ mass points of mass $m$ coupled by a mechanical spring to each other. Let the distance of the mass points from each other projected to the $x$-axis be $\Delta$ and introduce the transverse coordinates $q_{r}(t), r=1, \ldots, N$ of the mass points. In the limit $\Delta \rightarrow 0$ with $L$ fixed, the profile $q_{r}(t)$ asymptotes to the field $\phi(x, t)$. In this sense the field variable $x$ is the continuous label for infinitely many degrees of freedom.
We can now linearise the force between the mass points due to the spring. As a result of a simple exercise in classical mechanics the energy at leading order is found to be

$$
\begin{equation*}
E=\sum_{r=0}^{N}\left(\frac{1}{2} m\left(\frac{\mathrm{~d} q_{r}(t)}{\mathrm{d} t}\right)^{2}+k\left(q_{r}^{2}-q_{r} q_{r-1}\right)\right)+O\left(q^{3}\right), \quad k=\frac{T}{L} . \tag{1.4}
\end{equation*}
$$

In the continuum limit this becomes

$$
\begin{equation*}
E=\int_{0}^{L}\left[\frac{1}{2} \rho\left(\frac{\partial \phi(x, t)}{\partial t}\right)^{2}+\frac{1}{2} \rho c^{2}\left(\frac{\partial \phi(x, t)}{\partial x}\right)^{2}\right] \mathrm{d} x \tag{1.5}
\end{equation*}
$$

in terms of the mass density $\rho$ of the string and a suitably defined characteristic velocity $c$. Note that the second term indeed includes the nearest neighbour interaction because

$$
\begin{equation*}
\left(\frac{\partial \phi(x, t)}{\partial x}\right)^{2} \simeq\left(\lim _{\delta x \rightarrow 0} \frac{\phi(x+\delta x, t)-\phi(x, t))}{\delta x}\right)^{2} \tag{1.6}
\end{equation*}
$$

contains the off-diagonal terms $\phi(x+\delta x, t) \phi(x, t)$.
The nearest-neighbour interaction implies that the equation of motion for the mass points $q_{i}$ obey coupled linear differential equations. This feature persists in the continuum limit. To solve the dynamics it is essential that we are able to diagonalise the interaction in terms of the Fourier modes,

$$
\begin{align*}
\phi(x, t) & =\sum_{k=1}^{\infty} A_{k}(t) \sin \left(\frac{k \pi x}{L}\right),  \tag{1.7}\\
E & =\frac{L}{2} \sum_{k=1}^{\infty}\left(\frac{1}{2} \rho \dot{A}_{k}^{2}+\frac{1}{2} \rho \omega_{k}^{2} A^{2}\right),
\end{align*}
$$

where $\omega_{k}=k \pi c / L$. We are now dealing with a collection of infinitely many, decoupled harmonic oscillators $A_{k}(t)$.

In a final step, we quantise this collection of harmonic oscillators. According to Quantum Mechanics, each mode $A_{k}(t)$ can take energy values

$$
\begin{equation*}
E_{k}=\hbar \omega_{k}\left(n_{k}+1 / 2\right) \quad n_{k}=0,1,2, \ldots, \infty . \tag{1.8}
\end{equation*}
$$

The total energy is given by summing over the energy associated with all the modes, $E=\sum E_{k}$. A state of definite energy $E$ corresponds to mode numbers ( $n_{1}, n_{2}, \ldots, n_{\infty}$ ), where we think of $n_{r}$ as an excitation of the string or of the field $\phi$, i.e. as a quantum. In condensed matter physics, these quantised excitations in terms of harmonic modes are called quasi-particles, e.g. phonons for mechanical vibrations of a solid. Note that the above decoupling of the degrees of freedom rested on the quadratic form of the potential. Including higher terms will destroy this and induce interactions between modes.

The idea of Quantum Field Theory is to adapt this logic to particle physics and to describe a particle as the quantum of oscillation of an abstract field - just like in solid state physics we think of a quasi-particle as the vibrational excitation of a solid. The only difference is that the fields are now more abstract objects defined all over spacetime as opposed to concrete mechanical fields of the type above.
As a familiar example for a field we can think of the Maxwell field $A_{\mu}(x, t)$ in classical electrodynamics. A photon is the quantum excitation of this. It has spin 1 . Similarly we assign one field to each particle species, e.g. an electron is the elementary excitation of the electron field (Spin 1/2). We will interpret the sum over harmonic oscillator energies as an integral over possible energies for given momentum,

$$
\begin{equation*}
E=\sum_{k=1}^{\infty} \hbar \omega_{k}\left(n_{k}+1 / 2\right) \rightarrow \int \mathrm{d} p \hbar \omega_{p}\left(n_{p}+1 / 2\right) . \tag{1.9}
\end{equation*}
$$

A single particle with momentum $p$ corresponds to $n_{p}=1$ while all others vanish, but this is just a special example of a more multi-particle state with several $n_{p_{i}} \neq 0$. In particular, in agreement with the requirements of a multi-particle framework, at fixed $E$ transitions between various multi-particle states are in principle possible. Such transitions are induced by interactions corresponding to the higher order terms in the Hamiltonian that we have discarded so far. As a triumph this formalism also solves the problem of causality, as we will see.

### 1.2 Classical scalar field: Lagrangian formulation

We now formalise the outlined transition from a classical system with a finite number of degrees of freedom $q_{i}(t)$ to a classical field theory in terms of a scalar field $\phi(t, \vec{x}) \equiv \phi\left(x^{\mu}\right)$. In classical mechanics we start from an action

$$
\begin{equation*}
S=\int_{t_{1}}^{t_{2}} \mathrm{~d} t L\left(q_{i}(t), \dot{q}_{i}(t)\right) \text { with } L=\frac{1}{2} \sum_{i}\left(\dot{q}_{i}(t)\right)^{2}-V\left(q_{1},, q_{N}\right), \tag{1.10}
\end{equation*}
$$

where we have included the mass $m$ in the definition of $q_{i}(t)$. In a first step replace

$$
\begin{align*}
q_{i} & \rightarrow \phi\left(x^{\mu}\right) \equiv \phi(x),  \tag{1.11}\\
\dot{q}_{i}(t) & \rightarrow \frac{\partial \phi(x)}{\partial t} \tag{1.12}
\end{align*}
$$

thereby substituting the label $i=1, \ldots N$ by a continous coordinate $\vec{x} \equiv x^{i}$ with $i=1,2,3$. For the moment we consider a real scalar field i.e. $\phi(x)=\phi^{*}(x)$ which takes values in $\mathbb{R}$, i.e.

$$
\begin{equation*}
\phi: x^{\mu} \rightarrow \phi\left(x^{\mu}\right) \in \mathbb{R} . \tag{1.13}
\end{equation*}
$$

We will see that such a field describes spin-zero particles. Examples of scalar particles in nature are the Higgs boson or the inflaton, which cosmologists believe to be responsible for the exponential expansion of the universe during in inflation.

To set up the Lagrange function we first note that in a relativistic theory the partial time derivative can only appear as part of

$$
\begin{equation*}
\partial_{\mu} \phi(x) \equiv \frac{\partial}{\partial x^{\mu}} \phi(x) . \tag{1.14}
\end{equation*}
$$

Thus the Lagrange function can be written as

$$
\begin{equation*}
L=\int \mathrm{d}^{3} x \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right) \tag{1.15}
\end{equation*}
$$

where $\mathcal{L}$ is the Lagrange density. The action therefore is

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right) \tag{1.16}
\end{equation*}
$$

While, especially in condensed matter physics, also non-relativistic field theories are relevant, we focus on relativistic theories in this course.

Note furthermore that throughout this course we use conventions where

$$
\begin{equation*}
\hbar=c=1 . \tag{1.17}
\end{equation*}
$$

Then $\mathcal{L}$ has the dimension mass ${ }^{4}$, i.e. $[\mathcal{L}]=4$, since $[S]=0$ and $\left[\mathrm{d}^{4} x\right]=-4$.
The next goal is to find the Lagrangian: In a relativistic setting $\mathcal{L}$ can contain powers of $\phi$ and ${ }^{1}$ $\partial_{\mu} \phi \partial^{\mu} \phi \equiv \eta^{\mu \nu} \partial_{\mu} \phi \partial_{\nu} \phi$, which is the simplest scalar which can be built from $\partial_{\mu} \phi$. The action in this case is

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left[\frac{1}{2} \partial_{\mu} \phi \partial^{\mu} \phi-V(\phi)+O\left(\phi^{n}(\partial \phi)^{m}\right)\right], \tag{1.18}
\end{equation*}
$$

where

$$
\begin{equation*}
\frac{1}{2} \partial_{\mu} \phi \partial^{\mu} \phi=\frac{1}{2} \dot{\phi}^{2}-\frac{1}{2}(\nabla \phi)^{2} \tag{1.19}
\end{equation*}
$$

[^1]and the last type of terms consists of higher derivative terms with $m \geq 2$ or mixed terms with $n \geq 1$. Notice that

- the signature for the metric is, in our conventions, $(+,-,-,-)$, such that the sign in the action is indeed chosen correctly such that the kinetic term appears with a positive prefactor;
- $\phi$ has dimension 1 (mass ${ }^{1}$ ).

The potential $V(\phi(x))$ is in general a power series of the form

$$
\begin{equation*}
V(\phi(x))=a+b \phi(x)+c \phi^{2}(x)+d \phi^{3}(x)+\ldots \tag{1.20}
\end{equation*}
$$

We assume that the potential has a global minimum at $\phi(x)=\phi_{0}(x)$ such that

$$
\begin{equation*}
\left.\frac{\partial V(\phi)}{\partial \phi}\right|_{\phi=\phi_{0}}=0, \quad V\left(\phi_{0}\right)=V_{0} \tag{1.21}
\end{equation*}
$$

By a field redefinition we ensure that the minimum is at $\phi_{0}(x) \equiv 0$ and expand $V(\phi(x))$ around this minimum as

$$
\begin{equation*}
V(\phi(x))=V_{0}+\frac{1}{2} m^{2} \phi^{2}(x)+\boldsymbol{O}\left(\phi^{3}(x)\right) \tag{1.22}
\end{equation*}
$$

Here we used that the linear terms vanish at the extremum and the assumption that we are expanding around a minimum implies $m^{2}>0$. The constant $V_{0}$ is the classical contribution to the ground state or vacuum energy. Since in a theory without gravity absolute energies are not measurable, we set $V_{0}=0$ for the time being, but keep in mind that in principle $V_{0}$ is arbitrary. We will have considerably more to say about $V_{0}$ in the quantum theory in section (1.8).
Therefore the action becomes

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left[\frac{1}{2}(\partial \phi)^{2}-\frac{1}{2} m^{2} \phi^{2}+\ldots\right] \tag{1.23}
\end{equation*}
$$

We will find that $m^{2}$, the prefactor of the quadratic term, is related to the mass of the particles and that the omitted higher powers of $\phi$ as well as the terms $O\left(\phi^{n}(\partial \phi)^{m}\right)$ will give rise to interactions between these particles.
As an aside note that a negative value of $m^{2}$ signals that the extremum around which we are expanding the potential is a maximum rather than a minimum. Therefore $m^{2}<0$ signals a tachyonic instability: quantum fluctuations will destabilise the vacuum and cause the system to roll down its potential until it has settled in its true vacuum.
We will start by ignoring interaction terms and studying the action of the free real scalar field theory

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left[\frac{1}{2}(\partial \phi)^{2}-\frac{1}{2} m^{2} \phi^{2}\right], \phi=\phi^{*} \tag{1.24}
\end{equation*}
$$

The equations of motion are given by the Euler-Lagrange equations. As in classical mechanics we derive them by varying $S$ with respect to $\phi$ and $\partial_{\mu} \phi$ subject to $\left.\delta \phi\right|_{\text {boundary }}=\left.\delta \partial_{\mu} \phi\right|_{\text {boundary }}=0$. This
yields

$$
\begin{align*}
0 \stackrel{!}{=} \delta S & =\int \mathrm{d}^{4} x\left[\frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial \phi(x)} \delta \phi(x)+\frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial\left(\partial_{\mu} \phi(x)\right)} \delta \partial_{\mu} \phi(x)\right]  \tag{1.25}\\
& =\int \mathrm{d}^{4} x\left[\frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial \phi(x)} \delta \phi(x)+\frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial\left(\partial_{\mu} \phi(x)\right)} \partial_{\mu} \delta \phi(x)\right] .
\end{align*}
$$

Integrating by parts gives

$$
\begin{equation*}
\int \mathrm{d}^{4} x\left[\frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial \phi(x)}-\partial_{\mu} \frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial\left(\partial_{\mu} \phi(x)\right)}\right] \delta \phi(x)+\text { boundary terms. } \tag{1.26}
\end{equation*}
$$

Since the boundary terms vanish by assumption, therefore the integrand has to vanish for all variations $\delta \phi(x)$. This yields the Euler-Lagrange equations

$$
\begin{equation*}
\frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial \phi(x)}=\partial_{\mu} \frac{\partial \mathcal{L}\left(\phi(x), \partial_{\mu} \phi(x)\right)}{\partial\left(\partial_{\mu} \phi(x)\right)} . \tag{1.27}
\end{equation*}
$$

By inserting (1.24) into (1.27) we find the equations of motion for the free scalar field

$$
\begin{equation*}
\partial_{\mu} \frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)}=\partial_{\mu} \partial^{\mu} \phi=\frac{\partial \mathcal{L}}{\partial \phi}=-m^{2} \phi, \tag{1.28}
\end{equation*}
$$

i.e. the Klein-Gordon equation

$$
\begin{equation*}
\left(\partial^{2}+m^{2}\right) \phi(x)=0 \tag{1.29}
\end{equation*}
$$

Note that (1.29) is a relativistic wave equation and that it is solved by

$$
\begin{equation*}
e^{ \pm i p x} \text { with } p \equiv p^{\mu}=\left(p^{0}, \vec{p}\right) \tag{1.30}
\end{equation*}
$$

subject to the dispersion relation $-p^{2}+m^{2}=0$, i.e. $p^{0}= \pm \sqrt{\vec{p}^{2}+m^{2}}$. We now set

$$
\begin{equation*}
E_{\vec{p}}:=\sqrt{\vec{p}^{2}+m^{2}} \equiv E_{p} \tag{1.31}
\end{equation*}
$$

and write the general solution of (1.29) in the form

$$
\begin{equation*}
\phi(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(f(\vec{p}) e^{-i p x}+g(\vec{p}) e^{i p x}\right) \tag{1.32}
\end{equation*}
$$

where $p:=\left(E_{p}, \vec{p}\right)$ and $f^{*}(\vec{p})=g(\vec{p})$ for real $\phi$ and $p x=p \cdot x=p^{\mu} x_{\mu}$.

### 1.3 Noether's Theorem

A key role in Quantum Field Theory is played by symmetries. We consider a field theory with Lagrangian $\mathcal{L}\left(\phi, \partial_{\mu} \phi\right)$. A symmetry of the theory is then defined to be a field transformation by which $\mathcal{L}$ changes at most by a total derivative such that the action stays invariant. This ensures that the equations of motion are also invariant. Symmetries and conservation laws are related by Noether's Theorem ${ }^{2}$ :

[^2]Every continuous symmetry in the above sense gives rise to a Noether current $j^{\mu}(x)$ such that

$$
\begin{equation*}
\partial_{\mu} j^{\mu}(x)=0 \tag{1.33}
\end{equation*}
$$

upon use of the equations of motion ( $\equiv$ "on-shell").

This can be proven as follows:
For a continuous symmetry we can write infinitesimally:

$$
\begin{equation*}
\phi \rightarrow \phi+\epsilon \delta \phi+O\left(\epsilon^{2}\right) \text { with } \delta \phi=X\left(\phi, \partial_{\mu} \phi\right) \tag{1.34}
\end{equation*}
$$

Off-shell (i.e. without use of the equations of motion) we know that

$$
\begin{equation*}
\mathcal{L} \rightarrow \mathcal{L}+\epsilon \delta \mathcal{L}+O\left(\epsilon^{2}\right) \quad \text { with } \quad \delta \mathcal{L}=\partial_{\mu} F^{\mu} \tag{1.35}
\end{equation*}
$$

for some $F^{\mu}$. Now, under an arbitrary transformation $\phi \rightarrow \phi+\epsilon \delta \phi$, which is not necessarily a symmetry, $\delta \mathcal{L}$ is given by

$$
\begin{align*}
\delta \mathcal{L} & =\frac{\partial \mathcal{L}}{\partial \phi} \delta \phi+\frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)} \delta\left(\partial_{\mu} \phi\right) \\
& =\partial_{\mu}\left[\frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)} \delta \phi\right]+\left[\frac{\partial \mathcal{L}}{\partial \phi}-\partial_{\mu} \frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)}\right] \delta \phi \tag{1.36}
\end{align*}
$$

If $\delta \phi=X$ is a symmetry, then $\delta \mathcal{L}=\partial_{\mu} F^{\mu}$. Setting

$$
\begin{equation*}
j^{\mu}=\frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)} X-F^{\mu} \tag{1.37}
\end{equation*}
$$

we therefore have

$$
\begin{equation*}
\partial_{\mu} j^{\mu}=-\left(\frac{\partial \mathcal{L}}{\partial \phi}-\partial_{\mu} \frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)}\right) X \tag{1.38}
\end{equation*}
$$

off-shell. Note that the terms in brackets are just the Euler-Lagrange equation. Thus, if we use the equations of motion, i.e. on-shell, $\partial_{\mu} j^{\mu}=0$.

This immediately yields the following Lemma:
Every continuous symmetry whose associated Noether current satisfies $j^{i}(t, \vec{x}) \rightarrow 0$ sufficiently fast for $|\vec{x}| \rightarrow \infty$ gives rise to a conserved charge $Q$ with

$$
\begin{equation*}
\dot{Q}=0 \tag{1.39}
\end{equation*}
$$

Indeed if we take

$$
\begin{equation*}
Q=\int_{\mathbb{R}^{3}} \mathrm{~d}^{3} x j^{0}(t, \vec{x}) \tag{1.40}
\end{equation*}
$$

then the total time derivative of $Q$ is given by

$$
\begin{align*}
\dot{Q} & =\int_{\mathbb{R}^{3}} \mathrm{~d}^{3} x \frac{\partial}{\partial t} j^{0} \\
& =-\int_{\mathbb{R}^{3}} \mathrm{~d}^{3} x \partial_{i} j^{i}(t, \vec{x})=0 \tag{1.41}
\end{align*}
$$

by assumption of sufficiently fast fall-off of $j^{i}(t, \vec{x})$. We used that $\partial_{\mu} j^{\mu}=0$ in the first step.

The technical assumption $j^{i}(t, \vec{x}) \rightarrow 0$ for $|\vec{x}| \rightarrow \infty$ is really an assumption of 'sufficiently fast fall-off' of the fields at spatial infinity, which is typically satisfied. Note that in a finite volume $V=$ const., the quantity $Q_{V}=\int_{V} \mathrm{~d} V j^{0}(t, \vec{x})$ satisfies local charge conservation,

$$
\begin{equation*}
\dot{Q}_{V}=-\int_{V} \mathrm{~d} V \nabla \vec{j}=-\int_{\partial V} \vec{j} \cdot \mathrm{~d} \vec{s} . \tag{1.42}
\end{equation*}
$$

We now apply Noether's theorem to deduce the canonical energy-momentum tensor: Under a global spacetime transformation $x^{\mu} \rightarrow x^{\mu}+\epsilon^{\mu}$ a scalar field $\phi\left(x^{\mu}\right)$ transforms like

$$
\begin{equation*}
\phi\left(x^{\mu}\right) \rightarrow \phi\left(x^{\mu}-\epsilon^{\mu}\right)=\phi\left(x^{\mu}\right)-\epsilon^{v} \underbrace{\partial_{\gamma} \phi\left(x^{\mu}\right)}_{\equiv X_{v}(\phi)}+O\left(\epsilon^{2}\right) . \tag{1.43}
\end{equation*}
$$

Because $\mathcal{L}$ is a local function of $x$ it transforms as

$$
\begin{align*}
\mathcal{L} \rightarrow \mathcal{L}-\epsilon^{\nu} \partial_{\nu} \mathcal{L} & =\mathcal{L}-\eta^{\mu}{ }_{\nu} \epsilon^{\nu} \partial_{\mu} \mathcal{L}  \tag{1.44}\\
& =\mathcal{L}-\epsilon^{\nu} \partial_{\mu} \eta^{\mu} \mathcal{L} .
\end{align*}
$$

For each component $v$ we therefore have a conserved current $\left(j^{\mu}\right)_{v}$ given by

$$
\begin{equation*}
\left(j^{\mu}\right)_{v}=\frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)} \underbrace{\partial_{\nu} \phi}_{\equiv X_{v}}-\underbrace{\eta^{\mu}{ }_{v} \mathcal{L}}_{\equiv\left(F^{\mu}\right)_{v}} . \tag{1.45}
\end{equation*}
$$

With both indices up, we arrive at the canonical energy-momentum tensor

$$
\begin{equation*}
T^{\mu \nu}=\frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi\right)} \partial^{v} \phi-\eta^{\mu \nu} \mathcal{L} \text { with } \partial_{\mu} T^{\mu \nu}=0 \text { on-shell. } \tag{1.46}
\end{equation*}
$$

The conserved charges are the energy $E=\int d^{3} x T^{00}$ associated with time translation invariance and the spatial momentum $P^{i}=\int d^{3} x T^{0 i}$ associated with spatial translation invariance. We can combine them into the conserved 4-momentum

$$
\begin{equation*}
P^{v}=\int \mathrm{d}^{3} x T^{0 v} \tag{1.47}
\end{equation*}
$$

with the property $\dot{P}^{v}=0$.
Two comments are in order:

- In general, $T^{\mu v}$ may not be symmetric - especially in theories with spin. In such cases it can be useful to modify the energy-momentum tensor without affecting its conservedness or the associated conserved charges. Indeed we state as a fact that the Belinfante-Rosenfeld tensor

$$
\begin{equation*}
\Theta_{B R}^{\mu \nu}:=T^{\mu \nu}+\partial_{\rho} S^{\rho \mu \nu} \tag{1.48}
\end{equation*}
$$

can be defined in terms of a suitable $S^{\rho \mu \nu}=-S^{\mu \rho v}$ such that $\Theta_{B R}^{\mu \nu}$ is symmetric and obeys $\partial_{\mu} \Theta_{B R}^{\mu \nu}=0$.

- In General Relativity (GR), there exists yet another definition of the energy-momentum tensor: With the metric $\eta_{\mu \nu}$ replaced by $g_{\mu \nu}$ and

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x \sqrt{-g} \mathcal{L}^{\text {matter }}\left(g_{\mu v}, \phi, \partial \phi\right), \tag{1.49}
\end{equation*}
$$

where $g \equiv \operatorname{det} g$, one defines the Hilbert energy-momentum tensor:

$$
\begin{equation*}
\left(\Theta_{H}\right)^{\mu \nu}=-\frac{2}{\sqrt{-g}} \frac{\partial\left(\sqrt{-g} \mathcal{L}^{\text {matter }}\right)}{\partial g_{\mu v}} \tag{1.50}
\end{equation*}
$$

which is obviously symmetric and it the object that appears in the Einstein equations

$$
\begin{equation*}
R_{\mu \nu}+\frac{1}{2} R g_{\mu \nu}=8 \pi G\left(\Theta_{H}\right)_{\mu \nu} . \tag{1.51}
\end{equation*}
$$

In fact one can choose the Belinfante-Rosenfeld tensor such that it is equal to the Hilbert energymomentum tensor.

### 1.4 Quantisation in the Schrödinger Picture

Before quantising field theory let us briefly recap the transition from classical to quantum mechanics. We first switch from the Lagrange formulation to the canonical formalism of the classical theory. In classical mechanics the canonical momentum conjugate to $q_{i}(t)$ is

$$
\begin{equation*}
p_{i}(t)=\frac{\partial L}{\partial \dot{q}_{i}(t)} . \tag{1.52}
\end{equation*}
$$

The Hamiltonian is the Legendre transformation of the Lagrange function $L$

$$
\begin{equation*}
H=\sum_{i} p_{i}(t) \dot{q}_{i}(t)-L . \tag{1.53}
\end{equation*}
$$

To quantise in the Schrödinger picture we drop the time dependence of $q_{i}$ and $p_{i}$ and promote them to self-adjoint operators without any time dependence such that the fundamental commutation relation

$$
\begin{equation*}
\left[q_{i}, p_{j}\right]=i \delta_{i j} \tag{1.54}
\end{equation*}
$$

holds. Then all time dependence lies in the states.

This procedure is mimicked in a field theory by first defining clasically

$$
\begin{equation*}
\Pi(t, \vec{x}):=\frac{\partial \mathcal{L}}{\partial \dot{\phi}(t, \vec{x})} \tag{1.55}
\end{equation*}
$$

to be the conjugate momentum density. The Hamiltonian is

$$
\begin{equation*}
H=\int \mathrm{d}^{3} x \mathcal{H}=\int \mathrm{d}^{3} x[\Pi(t, \vec{x}) \dot{\phi}(t, \vec{x})-\mathcal{L}], \tag{1.56}
\end{equation*}
$$

where $\mathcal{H}$ is the Hamiltonian density. For the scalar field action (1.24) one finds

$$
\begin{equation*}
\Pi(t, \vec{x})=\dot{\phi}(t, \vec{x}) \tag{1.57}
\end{equation*}
$$

and therefore

$$
\begin{align*}
H & =\int \mathrm{d}^{3} x\left[\dot{\phi}^{2}(t, \vec{x})-\frac{1}{2}\left(\partial_{\mu} \phi\right)\left(\partial^{\mu} \phi\right)+\frac{1}{2} m^{2} \phi^{2}(t, \vec{x})\right] \\
& =\int \mathrm{d}^{3} x[\underbrace{\frac{1}{2} \dot{\phi}^{2}(t, \vec{x})}_{=\frac{1}{2} \Pi^{2}(t, \vec{x})}+\frac{1}{2}(\nabla \phi(t, \vec{x}))^{2}+\frac{1}{2} m^{2} \phi^{2}(t, \vec{x})] . \tag{1.58}
\end{align*}
$$

Note that as in classical mechanics one can define a Poisson bracket which induces a natural sympletic structure on phase space. In this formalism the Noether charges $Q$ are the generators of their underlying symmetries with the respect to the Poisson bracket (see Assignment 1 for details).

We now quantise in the Schrödinger picture. Therefore we drop the time-dependence of $\phi$ and $\Pi$ and promote them to Schrödinger-Picture operators $\phi^{(s)}(\vec{x})$ and $\Pi^{(s)}(\vec{x})$. For real scalar fields we get self-adjoint operators $\phi^{(s)}(\vec{x})=\left(\phi^{(s)}(\vec{x})\right)^{\dagger}$ with the canonical commutation relations (dropping ${ }^{(s)}$ from now on)

$$
\begin{equation*}
[\phi(\vec{x}), \Pi(\vec{y})]=i \delta^{(3)}(\vec{x}-\vec{y}), \quad[\phi(\vec{x}), \phi(\vec{y})]=0=[\Pi(\vec{x}), \Pi(\vec{y})] . \tag{1.59}
\end{equation*}
$$

### 1.5 Mode expansion

Our Hamiltonian (1.58) resembles the Hamiltonian describing a collection of harmonic oscillators, one at each point $\vec{x}$, but the term

$$
\begin{equation*}
(\nabla \phi)^{2} \approx\left(\frac{\phi(\vec{x}+\delta \vec{x})-\phi(\vec{x})}{|\delta \vec{x}|}\right)^{2} \tag{1.60}
\end{equation*}
$$

couples the degrees of freedom at $\vec{x}$ and $\vec{x}+\delta \vec{x}$. To arrive at a description in which the harmonic oscillators are decoupled, we must diagonalise the potential. Now, a basis of eigenfunctions with
respect to $\nabla$ is $e^{i \vec{p} \cdot \vec{x}}$. Thus the interaction will be diagonal in momentum space. With this motivation we Fourier-transform the fields as

$$
\begin{align*}
\phi(\vec{x}) & =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \tilde{\phi}(\vec{p}) e^{i \vec{p} \cdot \vec{x}} \\
\Pi(\vec{x}) & =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \tilde{\Pi}(\vec{p}) e^{i \vec{p} \cdot \vec{x}} \tag{1.61}
\end{align*}
$$

where $\tilde{\phi}^{\dagger}(\vec{p})=\tilde{\phi}(-\vec{p})$ ensures that $\phi(\vec{x})$ is self-adjoint. To compute $H$ in Fourier space we must insert these expressions into (1.58). First note that

$$
\begin{align*}
\frac{1}{2} \int \mathrm{~d}^{3} x(\nabla \phi(\vec{x}))^{2} & =\frac{1}{2} \int \mathrm{~d}^{3} x\left(\int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \nabla e^{i \vec{p} \cdot \vec{x}} \tilde{\phi}(\vec{p})\right)^{2}  \tag{1.62}\\
& =\frac{1}{2} \int \mathrm{~d}^{3} x \int \frac{\mathrm{~d}^{3} p \mathrm{~d}^{3} q}{(2 \pi)^{6}}(-\vec{p} \cdot \vec{q}) e^{i(\vec{p}+\vec{q}) \cdot \vec{x}} \tilde{\phi}(\vec{p}) \tilde{\phi}(\vec{q})
\end{align*}
$$

Thanks to the important equality

$$
\begin{equation*}
\int \mathrm{d}^{3} x e^{i(\vec{p}+\vec{q}) \cdot \vec{x}}=(2 \pi)^{3} \delta^{(3)}(\vec{p}+\vec{q}) \tag{1.63}
\end{equation*}
$$

the latter equation yields

$$
\begin{equation*}
\frac{1}{2} \int \mathrm{~d}^{3} x(\nabla \phi(\vec{x}))^{2}=\frac{1}{2} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \vec{p}^{2} \underbrace{\tilde{\phi}(\vec{p}) \tilde{\phi}(-\vec{p})}_{\equiv|\tilde{\phi}(\vec{p})|^{2}} \tag{1.64}
\end{equation*}
$$

and therefore with (1.58) altogether

$$
\begin{equation*}
H=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}}\left[\frac{1}{2}|\tilde{\Pi}(\vec{p})|^{2}+\frac{1}{2} \omega_{p}^{2}|\tilde{\phi}(\vec{p})|^{2}\right] \tag{1.65}
\end{equation*}
$$

where $\omega_{p}=\sqrt{\vec{p}^{2}+m^{2}}$. This is a collection of decoupled harmonic oscillators of frequency $\omega_{p}$ - and indeed the Hamiltonian is diagonal in momentum space.

The next step is to solve these oscillators in close analogy with the quantum mechanical treatment of a harmonic oscillator of frequency $\omega$ with Hamiltonian

$$
\begin{equation*}
H=\frac{1}{2} \Pi^{2}+\omega^{2} q^{2} \tag{1.66}
\end{equation*}
$$

The quantum mechanical definition of ladder operators $a, a^{\dagger}$ such that

$$
\begin{equation*}
q=\frac{1}{\sqrt{2 \omega}}\left(a+a^{\dagger}\right), \Pi=-\sqrt{\frac{\omega}{2}} i\left(a-a^{\dagger}\right) \tag{1.67}
\end{equation*}
$$

and with commutation relation $\left[a, a^{\dagger}\right]=1$ can be generalised to field theory as follows: By taking into account that

$$
\begin{equation*}
\tilde{\phi}(\vec{p})=\tilde{\phi}^{\dagger}(-\vec{p}), \quad \tilde{\Pi}(\vec{p})=\tilde{\Pi}^{\dagger}(-\vec{p}) \tag{1.68}
\end{equation*}
$$

we define the operators

$$
\begin{align*}
a(\vec{p}) & =\frac{1}{2}\left[\sqrt{2 \omega_{p}} \tilde{\phi}(\vec{p})+i \sqrt{\frac{2}{\omega_{p}}} \tilde{\Pi}(\vec{p})\right],  \tag{1.69}\\
a^{\dagger}(\vec{p}) & =\frac{1}{2}\left[\sqrt{2 \omega_{p}} \tilde{\phi}(-\vec{p})-i \sqrt{\frac{2}{\omega_{p}}} \tilde{\Pi}(-\vec{p})\right] .
\end{align*}
$$

Solving for $\tilde{\phi}(\vec{p})$ and $\tilde{\Pi}(\vec{p})$ and plugging into (1.61) yields

$$
\begin{align*}
& \phi(\vec{x})=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 \omega_{p}}}\left(a(\vec{p}) e^{i \vec{p} \cdot \vec{x}}+a^{\dagger}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right), \\
& \Pi(\vec{x})=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}}(-i) \sqrt{\frac{\omega_{p}}{2}}\left(a(\vec{p}) e^{i \vec{p} \cdot \vec{x}}-a^{\dagger}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right) . \tag{1.70}
\end{align*}
$$

From (1.61) and (1.59) we furthermore deduce the commutation relations

$$
\begin{equation*}
[\tilde{\phi}(\vec{p}), \tilde{\Pi}(\vec{q})]=\int \mathrm{d}^{3} x \mathrm{~d}^{3} y e^{-i \vec{p} \cdot \vec{x}} e^{-i \vec{q} \cdot \vec{y}} \underbrace{[\phi(\vec{x}), \Pi(\vec{y})]}_{=i \delta^{(3)}(\vec{x}-\vec{y})}=(2 \pi)^{3} i \delta^{(3)}(\vec{p}+\vec{q}), \tag{1.71}
\end{equation*}
$$

where again (1.63) was used, and

$$
\begin{equation*}
[\tilde{\phi}(\vec{p}), \tilde{\phi}(\vec{q})]=0=[\tilde{\Pi}(\vec{p}), \tilde{\Pi}(\vec{q})] . \tag{1.72}
\end{equation*}
$$

The ladder operators therefore obey the commutation relation

$$
\begin{array}{|l}
\hline\left[a(\vec{p}), a^{\dagger}(\vec{q})\right]=(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{q}),  \tag{1.73}\\
{\left[a^{\dagger}(\vec{p}), a^{\dagger}(\vec{q})\right]=0=[a(\vec{p}), a(\vec{q})]}
\end{array}
$$

The Hamiltonian (1.65) in mode expansion is

$$
\begin{align*}
H= & \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}}\left[\frac{1}{2}\left(i \sqrt{\frac{\omega_{p}}{2}}\right)^{2}\left(a(\vec{p})-a^{\dagger}(-\vec{p})\right)\left(a(-\vec{p})-a^{\dagger}(\vec{p})\right)\right.  \tag{1.74}\\
& \left.+\frac{\omega_{p}^{2}}{2} \frac{1}{2 \omega_{p}}\left(a(\vec{p})+a^{\dagger}(-\vec{p})\right)\left(a(-\vec{p})+a^{\dagger}(\vec{p})\right)\right] .
\end{align*}
$$

Only cross-terms survive due to the commutation relations:

$$
\begin{align*}
H & =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{\omega_{p}}{4}\left[a(\vec{p}) a^{\dagger}(\vec{p})+a^{\dagger}(-\vec{p}) a(-\vec{p})\right] \cdot 2 \\
& =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{\omega_{p}}{2}\left[a(\vec{p}) a^{\dagger}(\vec{p})+a^{\dagger}(-\vec{p}) a(-\vec{p})\right]  \tag{1.75}\\
& =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{\omega_{p}}{2}\left[a^{\dagger}(\vec{p}) a(\vec{p})+(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{p})+a^{\dagger}(-\vec{p}) a(-\vec{p})\right]
\end{align*}
$$

Therefore $H$ in its final form is given by

$$
\begin{equation*}
H=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \omega_{p} a^{\dagger}(\vec{p}) a(\vec{p})+\Delta_{H} \tag{1.76}
\end{equation*}
$$

where we renamed $-\vec{p} \rightarrow \vec{p}$ in the second term. The additional constant $\Delta_{H}$ is

$$
\begin{equation*}
\Delta_{H}=\frac{1}{2} \int \mathrm{~d}^{3} p \omega_{p} \delta^{(3)}(0) \tag{1.77}
\end{equation*}
$$

which is clearly divergent. An interpretation will be given momentarily. By explicit computation one finds that $H$ obeys the commutation relations

$$
\begin{align*}
& {[H, a(\vec{p})]=-\omega_{p} a(\vec{p}),}  \tag{1.78}\\
& {\left[H, a^{\dagger}(\vec{p})\right]=\omega_{p} a^{\dagger}(\vec{p}) .}
\end{align*}
$$

Similarly one computes the spatial momentum operator

$$
\begin{equation*}
P^{i}=\int \mathrm{d}^{3} x \dot{\phi}(\vec{x}) \partial^{i} \phi(\vec{x}) \tag{1.79}
\end{equation*}
$$

to be

$$
\begin{equation*}
P^{i}=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} p^{i} a^{\dagger}(\vec{p}) a(\vec{p})+\Delta_{p^{i}} \tag{1.80}
\end{equation*}
$$

with

$$
\begin{equation*}
\Delta_{p^{i}}=\frac{1}{2} \int \mathrm{~d}^{3} p p^{i} \delta^{(3)}(0) \equiv 0 \tag{1.81}
\end{equation*}
$$

We combine $H$ and $P$ into the 4-momentum operator

$$
\begin{equation*}
P^{\mu}=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} p^{\mu} a^{\dagger}(\vec{p}) a(\vec{p})+\Delta_{p^{\mu}} \tag{1.82}
\end{equation*}
$$

with $p^{\mu}=\left(p^{0}, \vec{p}\right)=\left(\omega_{p}, \vec{p}\right)$. It obeys the commutation relations

$$
\begin{align*}
& {\left[P^{\mu}, a^{\dagger}(\vec{p})\right]=p^{\mu} a^{\dagger}(\vec{p})}  \tag{1.83}\\
& {\left[P^{\mu}, a(\vec{p})\right]=-p^{\mu} a(\vec{p})}
\end{align*}
$$

### 1.6 The Fock space

We now find the Hilbert space on which the 4-moment operator $P^{\mu}$ acts. The logic is analogous to the considerations leading to the representation theory of the harmonic oscillator in Quantum Mechanics:

- Since $P^{\mu}$ is self-adjoint it has eigenstates with real eigenvalues. Let $\left|k^{\mu}\right\rangle$ be such an eigenstate with

$$
\begin{equation*}
P^{\mu}\left|k^{\mu}\right\rangle=k^{\mu}\left|k^{\mu}\right\rangle \tag{1.84}
\end{equation*}
$$

Then as a result of (1.83)

$$
\begin{align*}
P^{\mu} a^{\dagger}(\vec{q})\left|k^{\mu}\right\rangle & =a^{\dagger}(\vec{q}) P^{\mu}\left|k^{\mu}\right\rangle+q^{\mu} a^{\dagger}(\vec{q})\left|k^{\mu}\right\rangle  \tag{1.85}\\
& =\left(k^{\mu}+q^{\mu}\right) a^{\dagger}(\vec{q})\left|k^{\mu}\right\rangle
\end{align*}
$$

and similarly

$$
\begin{equation*}
P^{\mu} a(\vec{q})\left|k^{\mu}\right\rangle=\left(k^{\mu}-q^{\mu}\right) a(\vec{q})\left|k^{\mu}\right\rangle . \tag{1.86}
\end{equation*}
$$

This means that $a(\vec{q})$ and $a^{\dagger}(\vec{q})$ are indeed ladder operators which respectively subtract and add 4-momentum $q^{\mu}$ to or from $\left|k^{\mu}\right\rangle$.

- Next we observe that the Hamiltonian $H=P^{0}$ given by (1.76) is non-negative, i.e. $\langle\psi| H|\psi\rangle \geq$ $0 \forall$ states $|\psi\rangle$ because

$$
\begin{equation*}
\langle\psi| H|\psi\rangle=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \omega_{p}\langle\psi| a^{\dagger}(\vec{p}) a(\vec{p})|\psi\rangle+\Delta_{H}\langle\psi \mid \psi\rangle \geq 0 . \tag{1.87}
\end{equation*}
$$

Thus there exists a state $|0\rangle$ such that

$$
\begin{equation*}
a(\vec{q})|0\rangle=0 \forall \vec{q} . \tag{1.88}
\end{equation*}
$$

Otherwise successive action of $a(\vec{q})$ would lead to negative eigenvalues of $H .|0\rangle$ is called the vacuum of the theory. It has 4-momentum

$$
P^{\mu}|0\rangle=\Delta_{p^{\mu}}|0\rangle=\left\{\begin{array}{ll}
\Delta_{H}, & \mu=0  \tag{1.89}\\
0, & \mu=i
\end{array} .\right.
$$

- We interpret the divergent constant $\Delta_{H}$ given in (1.77) as the vacuum energy. A more thorough discussion of the significance of the divergence will be given later. For now we should note that in a theory without gravity absolute energy has no meaning. We can hus discard the additive constant $\Delta_{p^{\mu}}$ by defining

$$
\begin{equation*}
\tilde{P}^{\mu}:=P^{\mu}-\Delta_{p^{\mu}}=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} p^{\mu} a^{\dagger}(\vec{p}) a(\vec{p}), \tag{1.90}
\end{equation*}
$$

with $\tilde{P}^{\mu}|0\rangle=0$. From now on we only work with $\tilde{P}^{\mu}$ and drop the tilde.

- The state $a^{\dagger}(\vec{p})|0\rangle$ then has 4-momentum $p^{\mu}$,

$$
\begin{equation*}
P^{\mu} a^{\dagger}(\vec{p})|0\rangle=p^{\mu} a^{\dagger}(\vec{p})|0\rangle, \tag{1.91}
\end{equation*}
$$

with $p^{\mu}=\left(E_{p}, \vec{p}\right)$ and $E_{p}=\sqrt{\vec{p}^{2}+m^{2}}$. Since this is the relativistic dispersion relation for a single particle with mass $m$ we interpret $a^{\dagger}(\vec{p})|0\rangle$ as a 1-particle state with energy $E_{p}$ and momentum $\vec{p}$.

- More generally an N-particle state with energy $E=E_{p_{1}}+\ldots+E_{p_{N}}$ and momentum $\vec{p}=$ $\vec{p}_{1}+\ldots \vec{p}_{N}$ is given by

$$
\begin{equation*}
a^{\dagger}\left(\vec{p}_{1}\right) a^{\dagger}\left(\overrightarrow{p_{2}}\right) \ldots a^{\dagger}\left(\vec{p}_{N}\right)|0\rangle \tag{1.92}
\end{equation*}
$$

So much about the formalism. To get a better feeling for the objects we have introduced, let us recap what we have done: We have started with the assertion that spacetime - in our case $\mathbb{R}^{1,3}$ - is filled with the real scalar field $\phi(\vec{x})$, which we have taken to be a free field with Lagrangian $\mathcal{L}=\frac{1}{2}(\partial \phi)^{2}-\frac{1}{2} m^{2} \phi^{2}$. This field is interpreted as a field operator, i.e. in the Schrödinger picture at every space point $\vec{x}$ the object $\phi(\vec{x})$ represents a self-adjoint operator that acts on a Hilbert space. This Hilbert space possesses a state of lowest energy, the vacuum $|0\rangle$. The vacuum corresponds to the absence of any excitations of the field $\phi$ (at least on-shell). If one pumps energy $E_{p}$ and momentum $\vec{p}$ into some region of spacetime such that the relativistic dispersion relation $E_{p}=\sqrt{\vec{p}^{2}+m^{2}}$ holds, a particle $a^{\dagger}(\vec{p})|0\rangle$ is created as an excitation of $\phi(\vec{x})$. In particular, for the free theory the parameter $m$ in $\mathcal{L}$ is interpreted as the mass of such a particle. Since the underlying field $\phi(x)$ is a scalar field, the associated particle is called scalar particle. This realizes the shift of paradigm advertised at the very beginning of this course that the fundamental entity in Quantum Field Theory is not the particle, but rather the field:

The field $\phi(\vec{x})$ is the property of spacetime that in the presence of energy and momentum $\left(E_{p}, \vec{p}\right)$ a particle of energy $\left(E_{p}, \vec{p}\right)$ can be created.

In particular, this naturally gives rise to a multi-particle theory. The particles are just the excitations of the field and transitions with varying particle number can occur as long as the kinematics allows it. A first new, non-trivial conclusion we can draw from the formalism developed so far is the following special case of the spin-statistics theorem:

Scalar particles obey Bose statistics.
All we have to show that the N -particle wavefunction is symmetric under permutations. This follows immediately from the commutation relations, specifically the second line in (1.73):

$$
\begin{align*}
\left|\vec{p}_{1}, \ldots, \vec{p}_{i}, \ldots, \vec{p}_{j}, \ldots, \vec{p}_{n}\right\rangle & \simeq a^{\dagger}\left(\vec{p}_{1}\right) \ldots a^{\dagger}\left(\vec{p}_{i}\right) \ldots a^{\dagger}\left(\vec{p}_{j}\right) \ldots a^{\dagger}\left(\vec{p}_{N}\right)|0\rangle \\
& =a^{\dagger}\left(\vec{p}_{1}\right) \ldots a^{\dagger}\left(\vec{p}_{j}\right) \ldots a^{\dagger}\left(\vec{p}_{i}\right) \ldots a^{\dagger}\left(\vec{p}_{N}\right)|0\rangle  \tag{1.93}\\
& \simeq\left|\vec{p}_{1}, \ldots, \vec{p}_{j}, \ldots, \vec{p}_{i}, \ldots, \vec{p}_{n}\right\rangle,
\end{align*}
$$

where we have not fixed the normalization of the Fock state yet.

### 1.7 Some important technicalities

### 1.7.1 Normalisation

For reasons that will become clear momentarily, we choose to normalise the 1-particle momentum eigenstates as $|\vec{p}\rangle:=\sqrt{2 E_{p}} a^{\dagger}(\vec{p})|0\rangle$ and, more generally,

$$
\begin{equation*}
\left|\vec{p}_{1}, \ldots, \vec{p}_{n}\right\rangle:=\sqrt{2 E_{p_{1}} \cdot \ldots \cdot 2 E_{p_{N}}} a^{\dagger}\left(\vec{p}_{1}\right) \ldots a^{\dagger}\left(\vec{p}_{N}\right)|0\rangle . \tag{1.94}
\end{equation*}
$$

Then

$$
\begin{equation*}
\langle\vec{q} \mid \vec{p}\rangle=\sqrt{2 E_{p}} \sqrt{2 E_{q}}\langle 0| a(\vec{q}) a^{\dagger}(\vec{p})|0\rangle . \tag{1.95}
\end{equation*}
$$

To compute the inner product of two such states we use an important trick: Move all $a$ 's to the right and all $a^{\dagger}$ to the left with the help of

$$
\begin{equation*}
a(\vec{q}) a^{\dagger}(\vec{p})=a^{\dagger}(\vec{q}) a(\vec{p})+(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{q}) . \tag{1.96}
\end{equation*}
$$

This gives rise to terms of the form $a(\vec{q})|0\rangle=0=\langle 0| a^{\dagger}(\vec{p})$. Therefore

$$
\begin{equation*}
\langle\vec{q} \mid \vec{p}\rangle=(2 \pi)^{3} 2 E_{p} \delta^{(3)}(\vec{p}-\vec{q}) . \tag{1.97}
\end{equation*}
$$

Note that, as in Quantum Mechanics, momentum eigenstates are not strictly normalisable due to the appearance of the delta-distribution, but we can form normalisable states as wavepackets

$$
\begin{equation*}
|f\rangle=\int \mathrm{d}^{3} p f(\vec{p})|\vec{p}\rangle . \tag{1.98}
\end{equation*}
$$

### 1.7.2 The identity

With the above normalisation the identity operator on the 1-particle Hilbert space is

$$
\begin{equation*}
\mathbb{1}_{1-\text { particle }}=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}|\vec{p}\rangle\langle\vec{p}| . \tag{1.99}
\end{equation*}
$$

One should notice that $\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}$ is a Lorentz-invariant measure. This, in turn, is part of the motivation for the normalisation of the 1-particle momentum eigentstates. To see this we rewrite the measure as

$$
\begin{align*}
\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}} & =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{4}} \frac{2 \pi}{2 E_{p}} \\
& =\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} 2 \pi \delta\left(p^{2}-m^{2}\right) \Theta\left(p^{0}\right), \tag{1.100}
\end{align*}
$$

where we used that $\delta(a x)=\frac{1}{a} \delta(x)$ and that

$$
\begin{equation*}
\delta\left(p^{2}-m^{2}\right)=\delta\left(\left(p^{0}-E_{p}\right)\left(p^{0}+E_{p}\right)\right) \tag{1.101}
\end{equation*}
$$

in the last step. (1.100) is manifestly Lorentz-invariant: First, $\mathrm{d}^{4} p \rightarrow \operatorname{det}(\Lambda) \mathrm{d}^{4} p$ under a Lorentztransformation, and since the determinat of a Lorentz-transformation is 1 , it follows that $\mathrm{d}^{4} p$ is Lorentz-invariant. Moreover the sign of $p^{0}$ is unchanged under a Lorentz transformation.

### 1.7.3 Position-space representation

In Quantum Mechanics, the position eigenstate is related by a Fourier transformation to the momentum eigenstates, $|x\rangle=\int \frac{\mathrm{d} p}{2 \pi} e^{-i p x}|p\rangle$ with $\langle x \mid p\rangle=e^{i p x}$. Due to our normalisation

$$
\begin{equation*}
|\vec{p}\rangle=\sqrt{2 E_{p}} a^{\dagger}(\vec{p})|0\rangle \tag{1.102}
\end{equation*}
$$

the correct expression for $|\vec{x}\rangle$ in QFT is

$$
\begin{equation*}
|\vec{x}\rangle=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}} e^{-i \vec{p} \cdot \vec{x}}|\vec{p}\rangle \tag{1.103}
\end{equation*}
$$

because then

$$
\begin{equation*}
\langle\vec{x} \mid \vec{p}\rangle=\int \frac{\mathrm{d}^{3} q}{(2 \pi)^{3}} \frac{1}{2 E_{q}} e^{i \vec{q} \cdot \vec{x}}\langle\vec{q} \mid \vec{p}\rangle=e^{i \vec{p} \cdot \vec{x}}, \tag{1.104}
\end{equation*}
$$

where we used (1.97). Note that

$$
\begin{equation*}
|\vec{x}\rangle=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(e^{-i \vec{p} \cdot \vec{x}} a^{\dagger}(\vec{p})+e^{i \vec{p} \cdot \vec{x}} a(\vec{p})\right)|0\rangle=\phi(\vec{x})|0\rangle . \tag{1.105}
\end{equation*}
$$

In other words, the field operator $\phi(x)$ acting on the vacuum $|0\rangle$ creates a 1-particle position eigenstate.

### 1.8 On the vacuum energy

We had seen that originally

$$
\begin{equation*}
H=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \omega_{p} a^{\dagger}(\vec{p}) a(\vec{p})+\Delta_{H}, \quad \Delta_{H}=\frac{1}{2} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \omega_{p}(2 \pi)^{3} \delta(0) \tag{1.106}
\end{equation*}
$$

where $\Delta_{H}$ is the vacuum energy $E_{0}$ such that $H|0\rangle=\Delta_{H}|0\rangle \equiv E_{0}|0\rangle$. $E_{0}$ is the first example of a divergent quantity in QFT. In fact, it realises the two characteristic sources of a possible divergence in QFT:

- The divergent factor $(2 \pi)^{3} \delta^{(3)}(0)$ is interpreted as follows: We know that

$$
\begin{equation*}
\int_{\mathbb{R}^{3}} \mathrm{~d}^{3} x e^{i \vec{p} \cdot \vec{x}}=(2 \pi)^{3} \delta^{(3)}(\vec{p}), \tag{1.107}
\end{equation*}
$$

so formally the volume of $\mathbb{R}^{3}$ is given by

$$
\begin{equation*}
V_{\mathbb{R}^{3}}=\int_{\mathbb{R}^{3}} \mathrm{~d}^{3} x=(2 \pi)^{3} \delta^{(3)}(0) \tag{1.108}
\end{equation*}
$$

The divergence of $\delta^{(3)}(0)$ is rooted in the fact that the volume of $\mathbb{R}^{3}$ is infinite, and the corresponding divergent factor in $E_{0}$ arises because we are computing an energy in an infinite volume. This divergent factor thus results from the long-distance (i.e. small energy) behaviour of the theory and is an example of an infra-red (IR) divergence. Generally in QFT, IR divergences signal that we are either making a mistake or ask an unphysical question. In our case, the mistake is to consider the theory in a strictly infinite volume, which is of course unphysical. One can regularise the IR divergence by instead considering the theory in a given, but finite volume. What is free of the IR divergence is in particular the vacuum energy density

$$
\begin{equation*}
\epsilon_{0}=\frac{E_{0}}{V_{\mathbb{R}^{3}}}=\frac{1}{2} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \omega_{p} \tag{1.109}
\end{equation*}
$$

- Nonetheless, even $\epsilon_{0}$ remains divergent because

$$
\begin{equation*}
\epsilon_{0}=\frac{E_{0}}{V_{\mathbb{R}^{3}}}=\int \frac{\mathrm{d}^{3} p}{2(2 \pi)^{3}} \sqrt{\vec{p}^{2}+m^{2}}=\frac{1}{2} \frac{4 \pi}{(2 \pi)^{3}} \int_{0}^{\infty} \mathrm{d} p p^{2} \sqrt{p^{2}+m^{2}}, \tag{1.110}
\end{equation*}
$$

which goes to infinity due to the integration over all momenta up to $p \rightarrow \infty$. This is an ultraviolet (UV) divergence. The underlying reason for this (and all other UV divergences in QFT) is the breakdown of the theory at high energies (equivalently at short distances) - or at least a breakdown of our treatment of the theory.

To understand this last point it is beneficial to revisit the mechanical model of the field $\phi(t, x)$ as an excitation of a mechanical string as introduced in section 1.1. Recall that the field $\phi(t, x)$ describes the transverse position of the string in the continuum limit of vanishing distance $\Delta$ between the individual mass points at position $q_{i}(t)$ which were thought of as connected by an elastic spring. However, in reality the string is made of atoms of finite, typical size $R$. A continuous string profile $\phi(t, x)$ is therefore not an adequate description at distances $\Delta \leq R$ or equivalently at energies $E \geq \Lambda \simeq 1 / R$ resolving such small distances. Rather, if we want to describe processes at energies $E \geq \Lambda$, the continuous field theory $\phi(t, x)$ is to be replaced by the more fundamental, microscopic theory of atoms in a lattice. In this sense the field $\phi(t, x)$ gives merely an effective description of the string valid at energies $E \leq \Lambda$. Extrapolation of the theory beyond such energies is doubtful and can give rise to infinities - the UV divergences.
In a modern approach to QFT, this reasoning is believed to hold also for the more abstract relativistic fields we are considering in this course. According to this logic, QFT is really an effective theory that eventually must be replaced at high energies by a more fundamental theory. A necessary condition for such a fundamental theory to describe the microscopic degrees of freedom correctly is that it must be free of pathologies of all sort and in particular be UV finite. ${ }^{3}$ At the very least, gravitational degrees of freedom become important in the UV region and are expected to change the qualitative behaviour of the theory at energies around the Planck scale $M_{P} \simeq 10^{19} \mathrm{GeV}$.
Despite these limitations, in a 'good' QFT the UV divergences can be removed - for all practical purposes - by the powerful machinery of regularisation and renormalisation. We will study this procedure in great detail later in the course, but let us take this opportunity to very briefly sketch the logic for the example of the vacuum energy density:

- The first key observation is that in the classical Lagrangian we can have a constant term $V_{0}$ of dimension mass ${ }^{4}$,

$$
\begin{equation*}
\mathcal{L}=\frac{1}{2}(\partial \phi)^{2}-\frac{1}{2} m^{2} \phi^{2}-V_{0}, \quad V_{0}=\left.V\right|_{\min }, \tag{1.111}
\end{equation*}
$$

corresponding to the value of the classical potential at the minimum around which we expand. We had set $V_{0} \equiv 0$ in our analysis because we had argued that such an overall energy offset is not measurable in a theory without gravity. However, as we have seen the vacuum energy

[^3]density really consists of two pieces - the classical offset $V_{0}$ and the quantum piece $\Delta_{H}$. Thus, let us keep $V_{0}$ for the moment and derive the Hamiltonian
\[

$$
\begin{equation*}
H=\int \mathrm{d}^{3} x\left(\frac{1}{2} \Pi^{2}+\frac{1}{2} m^{2} \phi^{2}+\frac{1}{2}(\nabla \phi)^{2}\right)+\int \mathrm{d}^{3} x V_{0} . \tag{1.112}
\end{equation*}
$$

\]

- Next, we regularise the theory: Introduce a cutoff-scale $\Lambda$ and for the time being only allow for energies $E \leq \Lambda$. If we quantise the theory with such a cutoff at play, the overall vacuum energy is now

$$
\begin{equation*}
H|0\rangle=V_{\mathbb{R}^{3}}\left(\epsilon_{0}(\Lambda)+V_{0}\right)|0\rangle \tag{1.113}
\end{equation*}
$$

with

$$
\begin{equation*}
\epsilon_{0}(\Lambda)=\frac{1}{(2 \pi)^{2}} \int_{0}^{\Lambda} \mathrm{d} p p^{2} \sqrt{p^{2}+m^{2}} \tag{1.114}
\end{equation*}
$$

Note that the momentum integral only runs up to the cutoff $\Lambda$ in the regularised theory.

- Since $V_{0}$ is just a parameter, we can set

$$
\begin{equation*}
V_{0}=V_{0}(\Lambda)=-\epsilon_{0}(\Lambda)+\chi, \quad \chi \text { finite as } \Lambda \rightarrow \infty . \tag{1.115}
\end{equation*}
$$

With this choice

$$
\begin{equation*}
H|0\rangle=V_{\mathbb{R}^{3}} \chi|0\rangle \tag{1.116}
\end{equation*}
$$

independently of $\Lambda$ !This way we absorb the divergence into a cutoff-dependent counterterm $V_{0}(\Lambda)$ in the action such that the total vacuum energy density is finite. This step is called renormalisation. Crucially, note that the finite piece $\chi$ is completely arbitrary a priori. It must be determined experimentally by measuring a certain observable - in this case the vacuum energy (which is a meaningful observable only in the presence of gravity - see below).

- Finally, we can remove the cutoff by taking $\Lambda \rightarrow \infty$.

To summarize, we define the quantum theory as the result of quantising the classical Lagrangian $\mathcal{L}=\frac{1}{2}(\partial \phi)^{2}-\frac{1}{2} m^{2} \phi^{2}-V_{0}(\Lambda)$ with $V_{0}(\Lambda)=-\epsilon_{0}(\Lambda)+\chi$ and taking the limit $\Lambda \rightarrow \infty$ at the very end. Since $\Lambda$ appears only in the classical - or bare - Lagrangian, but in no observable, physical quantity at any of the intermediate stages, we can safely remove it at the end. In this sense, the theory is practically defined up to all energies.
This way to deal with UV divergences comes at a prize: We lose the prediction of one observable per type of UV divergence as a result of the inherent arbitrariness of the renormalisation step. In fact, above we have chosen $V_{0}(\Lambda)$ such that $H|0\rangle=V_{\mathbb{R}^{3}} \chi|0\rangle$. It is only in the absence of gravity that the vacuum energy is unobservable and thus $\chi$ is irrelevant. More generally, the value of the vacuum energy must now be taken from experiment $-\chi$ is an input parameter of the theory rather than a prediction. While we have succeeded in removing the divergence by renormalising the original Lagrangian,
the actual value of the physical observable associated with the divergence - here the vacuum energy density - must be taken as an input parameter from experiment or from other considerations. In a renormalisable QFT it is sufficient to do this for a finite number of terms in the Lagrangian so that once the associated observables are specified, predictive power is maintained for the computation of all subsequent observables.

## The Cosmological Constant in gravity

In gravity, the vacuum energy density is observable because it gravitates and we have to carry the vacuum energy with the field equations

$$
\begin{equation*}
R_{\mu \nu}-\frac{1}{2} R g_{\mu \nu}=-8 \pi G T_{\mu \nu}+\left(\epsilon_{0}+V_{0}\right) g_{\mu \nu} . \tag{1.117}
\end{equation*}
$$

This form of the Einstein field equations leads to accelerated expansion of the Universe. Indeed, observations indicate that the universe expands in an accelerated fashion, and the simplest - albeit not the only possible - explanation would be to identify the underlying 'Dark energy' with the vacuum energy. Observationally, this would then point to $\epsilon_{0}+V_{0}=\chi=\left(10^{-3} \mathrm{eV}\right)^{4}$. Of course in our QFT approach it is impossible to explain such a value of the vacuum energy because as a side-effect of renormalisation we gave up on predicting it. In this respect, Quantum Field Theory remains an effective description. In a fundamental, UV finite theory this would be different: There the net vacuum energy would be the difference of a finite piece $V_{0}$ and a finite quantum contribution $\epsilon_{0}$, both of which would be computable from first principles (if the theory is truly fundamental). Thus the two quantities should almost cancel each other. There is a problem, though: On dimensional grounds one expects both $V_{0}$ and $\epsilon_{0}$ to be of the order of the fundamental scale in the theory, which in a theory of quantum gravity is the Planck scale $M_{P}=10^{19} \mathrm{GeV}$. The expected value for the vacuum energy density is thus $M_{P}^{4}$, which differs by the observed value by about 122 orders of magnitude difference. Thus, the difference of $V_{0}$ and $\epsilon_{0}$ must be by 122 orders of magnitude smaller than both individual numbers. Such a behaviour is considered immense fine-tuning and thus unnatural. The famous Cosmological Constant Problem is therefore the puzzle of why the observed value is so small. ${ }^{4}$

### 1.9 The complex scalar field

We now extend the formalism developed so far to the theory of a complex scalar field, which classically no longer satisfies $\phi(x)=\phi^{*}(x)$. A convenient way to describe a complex scalar field of mass $m$ is to note that its real and imaginary part can be viewed as independent real scalar fields $\phi_{1}$ and $\phi_{2}$ of mass $m$, i.e. we can write

$$
\begin{equation*}
\phi(x)=\frac{1}{\sqrt{2}}\left(\phi_{1}(x)+i \phi_{2}(x)\right) . \tag{1.118}
\end{equation*}
$$

The real fields $\phi_{1}$ and $\phi_{2}$ are canonically normalised if we take as Lagrangian for the complex field

$$
\begin{equation*}
\mathcal{L}=\partial_{\mu} \phi^{\dagger}(x) \partial^{\mu} \phi(x)-m^{2} \phi^{\dagger}(x) \phi(x) . \tag{1.119}
\end{equation*}
$$

[^4]It is then a simple matter to repeat the programme of quantisation, e.g. by quantizing $\phi_{1}$ and $\phi_{2}$ as before and rewriting everything in terms of the complex field $\phi(x)$. At the end of this rewriting we can forget about $\phi_{1}$ and $\phi_{2}$ and simply describe the theory in terms of the complex field $\phi(x)$. The details of this exercise will be provided in the tutorials so that we can be brief here and merely summarise the main formulae.

- The fields $\phi(x)$ and $\phi^{\dagger}(x)$ describe independent degrees of freedom with respective conjugate momenta

$$
\begin{align*}
\Pi(t, \vec{x}) & =\frac{\partial \mathcal{L}}{\partial \dot{\phi}(t, \vec{x})}=\dot{\phi}^{\dagger}(t, \vec{x}) \\
\Pi^{\dagger}(t, \vec{x}) & =\frac{\partial \mathcal{L}}{\partial \dot{\phi}^{\dagger}(t, \vec{x})}=\dot{\phi}(t, \vec{x}) \tag{1.120}
\end{align*}
$$

The Hamiltonian $H$ is

$$
\begin{align*}
H & =\int \mathrm{d}^{3} x\left(\Pi^{\dagger}(\vec{x}) \dot{\phi}^{\dagger}(\vec{x})+\Pi(\vec{x}) \dot{\phi}(\vec{x})-\mathcal{L}\right) \\
& =\int \mathrm{d}^{3} x\left(\Pi^{\dagger}(\vec{x}) \Pi(\vec{x})+\nabla \phi^{\dagger}(\vec{x}) \nabla \phi(\vec{x})+m^{2} \phi^{\dagger}(\vec{x}) \phi(\vec{x})\right) \tag{1.121}
\end{align*}
$$

- These fields are promoted to Schrödinger picture operators with non-vanishing commutators

$$
\begin{equation*}
[\phi(\vec{x}), \Pi(\vec{y})]=i \delta^{(3)}(\vec{x}-\vec{y})=\left[\phi^{\dagger}(\vec{x}), \Pi^{\dagger}(\vec{y})\right] \tag{1.122}
\end{equation*}
$$

and all other commutators vanishing.

- The mode expansion is conveniently written as

$$
\begin{align*}
\phi(\vec{x}) & =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(a(\vec{p}) e^{i \vec{p} \cdot \vec{x}}+b^{\dagger}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right)  \tag{1.123}\\
\phi^{\dagger}(\vec{x}) & =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(b(\vec{p}) e^{i \vec{p} \cdot \vec{x}}+a^{\dagger}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right),
\end{align*}
$$

where the mode operators $a(\vec{p})$ and $b(\vec{p})$ are independent and $a^{\dagger}(\vec{p})$ and $b^{\dagger}(\vec{p})$ describe the respective conjugate operators. A quick way to arrive at this form of the expansion is to plug the mode expansion of the real fields $\phi_{1}$ and $\phi_{2}$ into (1.118). This identifies

$$
\begin{align*}
a & =\frac{1}{\sqrt{2}}\left(a_{1}+i a_{2}\right) \\
b^{\dagger} & =\frac{1}{\sqrt{2}}\left(a_{1}^{\dagger}+i a_{2}^{\dagger}\right) \tag{1.124}
\end{align*}
$$

In particular this implies that

$$
\begin{equation*}
\left[a(\vec{p}), a^{\dagger}(\vec{q})\right]=(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{q})=\left[b(\vec{p}), b^{\dagger}(\vec{q})\right] \tag{1.125}
\end{equation*}
$$

while all other commutators vanish.

- The mode expansion of the 4-momentum operator is

$$
\begin{equation*}
P^{\mu}=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} p^{\mu}\left(a^{\dagger}(\vec{p}) a(\vec{p})+b^{\dagger}(\vec{p}) b(\vec{p})\right) \tag{1.126}
\end{equation*}
$$

Crucially, one can establish that it has 2 types of momentum eigenstates

$$
\begin{equation*}
a^{\dagger}(\vec{p})|0\rangle \text { and } b^{\dagger}(\vec{p})|0\rangle \tag{1.127}
\end{equation*}
$$

both of energy $E_{p}=\sqrt{\vec{p}^{2}+m^{2}}$. I.e. both states have mass $m$, but they differ in their $U(1)$ charge as we will see now:

- The Lagrange density is invariant under the global continuous $U(1)$ symmetry

$$
\begin{equation*}
\phi(x) \rightarrow e^{i \alpha} \phi(x) \tag{1.128}
\end{equation*}
$$

where $\alpha$ is a constant in $\mathbb{R}$. Recall that the unitary group $U(N)$ is the group of complex $N \times N$ matrices $A$ satisfying $A^{\dagger}=A^{-1}$. The dimension of this group is $N^{2}$. In particular $e^{i \alpha} \in U(1)$. According to Noether's theorem there exists a conserved current (see Ass. 3)

$$
\begin{equation*}
j^{\mu}=-i\left(\phi^{\dagger} \partial^{\mu} \phi-\partial^{\mu} \phi^{\dagger} \phi\right) \tag{1.129}
\end{equation*}
$$

and charge

$$
\begin{equation*}
Q=\int \mathrm{d}^{3} x j^{0}=-\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}}\left(a^{\dagger}(\vec{p}) a(\vec{p})-b^{\dagger}(\vec{p}) b(\vec{p}) .\right) \tag{1.130}
\end{equation*}
$$

This Noether charge acts on a particle with momentum $\vec{p}$ as follows:

$$
\begin{align*}
& Q a^{\dagger}(\vec{p})|0\rangle=-a^{\dagger}(\vec{p})|0\rangle: \text { charge }-1 \\
& Q b^{\dagger}(\vec{p})|0\rangle=+b^{\dagger}(\vec{p})|0\rangle: \text { charge }+1 \tag{1.131}
\end{align*}
$$

We interpret $a^{\dagger}(\vec{p})|0\rangle$ as a particle of mass $m$ and charge -1 and $b^{\dagger}(\vec{p})|0\rangle$ as a particle with the same mass, but positive charge, i.e. as its anti-particle. For the real field the particle is its own anti-particle. Note that the term 'charge $\pm 1$ ' so far refers simply to the eigenvalue of the Noether charge operator $Q$ associated with the global $U(1)$ symmetry of the theory. That this abstract charge really coincides with what we usually call charge in physics - i.e. that it describes the coupling to a Maxwell type field - will be confirmed later when we study Quantum Electrodynamics.

### 1.10 Quantisation in the Heisenberg picture

So far all field operators have been defined in the Schrödinger picture, in which the time-dependence is carried entirely by the states on which these operators act. From Quantum Mechanics we recall
that alternatively quantum operators can be described in Heisenberg picture (HP), where the timedependence is carried by the operators $A^{(H)}(t)$ and not the states. The HP operator $A^{(H)}(t)$ is defined as

$$
\begin{equation*}
A^{(H)}(t)=e^{i H^{(S)}\left(t-t_{0}\right)} A^{(S)} e^{-i H^{(S)}\left(t-t_{0}\right)} \tag{1.132}
\end{equation*}
$$

where $A^{(S)}$ is the corresponding Schrödinger picture operator and $H^{(S)}$ is the Schrödinger picture Hamilton operator. At the time $t_{0}$ the Heisenberg operator and the Schrödinger operator coincide, i.e. $A^{(H)}\left(t_{0}\right)=A^{(S)}$. We will set $t_{0} \equiv 0$ from now on. The definition (1.132) has the following implications:

- $H^{(H)}(t)=H^{(S)} \forall t$.
- The time evolution of the Heisenberg picture operators is governed by the equation of motion

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} A^{(H)}(t)=i\left[H, A^{(H)}(t)\right] \tag{1.133}
\end{equation*}
$$

- The Schrödinger picture commutators translate into equal time commutators, e.g.

$$
\begin{equation*}
\left[q_{i}^{(H)}(t), p_{j}^{(H)}(t)\right]=i \delta_{i j} \tag{1.134}
\end{equation*}
$$

In field theory we similarly define the Heisenberg fields via

$$
\begin{align*}
\phi^{(H)}(t, \vec{x}) & \equiv \phi(x)=e^{i H^{(S)} t} \phi^{(S)}(\vec{x}) e^{-i H^{(S)} t} \\
\Pi^{(H)}(t, \vec{x}) & \equiv \Pi(x)=e^{i H^{(S)} t} \Pi^{(S)}(\vec{x}) e^{-i H^{(S)} t}  \tag{1.135}\\
\mathcal{H}^{(H)}(t, \vec{x}) & \equiv \mathcal{H}(x)=e^{i H^{(S)}} \mathcal{H}^{(S)}(\vec{x}) e^{-i H^{(S)} t}
\end{align*}
$$

These obey equal-time canonical commutation relations

$$
\begin{align*}
{[\phi(t, \vec{x}), \Pi(t, \vec{y})] } & =i \delta^{(3)}(\vec{x}-\vec{y})  \tag{1.136}\\
{[\phi(t, \vec{x}), \phi(t, \vec{y})] } & =0=[\Pi(t, \vec{x}), \Pi(t, \vec{y})]
\end{align*}
$$

The Heisenberg equation of motion for $\phi(x)$ reads

$$
\begin{equation*}
\frac{\partial}{\partial t} \phi(t, \vec{x})=i[H, \phi(t, \vec{x})]=i \int \mathrm{~d}^{3} y[\mathcal{H}(t, \vec{y}), \phi(t, \vec{x})] \tag{1.137}
\end{equation*}
$$

with

$$
\begin{equation*}
\mathcal{H}(t, \vec{y})=\frac{1}{2} \Pi^{2}(t, \vec{y})+\frac{1}{2}(\nabla \phi(t, \vec{y}))^{2}+\frac{1}{2} m^{2} \phi^{2}(t, \vec{y}) \tag{1.138}
\end{equation*}
$$

for a real scalar field. In the latter equation we used that $H$ is time-independent, i.e. we can evaluate the commutator at arbitrary times and thus choose equal time with $\phi(t, \vec{x})$ so that we can exploit the equal-time commutation relations.
In evaluating (1.137) we observe that the only non-zero term comes from

$$
\begin{equation*}
\frac{1}{2}\left[\Pi^{2}(t, \vec{y}), \phi(t, \vec{x})\right]=(-i) \Pi(t, \vec{y}) \delta^{(3)}(\vec{x}-\vec{y}) \tag{1.139}
\end{equation*}
$$

where used the standard relation $[A, B C]=[A, B] C+B[A, C]$ together with (1.136). This gives

$$
\begin{equation*}
\frac{\partial}{\partial t} \phi(t, \vec{x})=\Pi(t, \vec{x}) \tag{1.140}
\end{equation*}
$$

as expected. One can similarly show that

$$
\begin{align*}
\frac{\partial}{\partial t} \Pi(t, \vec{x}) & =i[H, \Pi(t, \vec{y})]  \tag{1.141}\\
& =\nabla^{2} \phi(t, \vec{x})-m^{2} \phi(t, \vec{x}) .
\end{align*}
$$

Therefore altogether we can establish that the Klein-Gordon equation

$$
\begin{equation*}
\left(\partial^{2}+m^{2}\right) \phi(x)=0 \tag{1.142}
\end{equation*}
$$

holds as an operator equation at the quantum level.
The covariantisation of (1.137) is

$$
\begin{equation*}
\partial^{\mu} \phi(x)=i\left[P^{\mu}, \phi(x)\right] \tag{1.143}
\end{equation*}
$$

with $P^{i}=\int \mathrm{d}^{3} y \Pi(y) \partial^{i} \phi(y)$. Indeed this equation can be explicitly confirmed by evaluating the commutator $\left[P^{i}, \phi(x)\right]$. As a consequence we will check, on sheet 3 , that

$$
\begin{equation*}
\phi\left(x^{\mu}+a^{\mu}\right)=e^{i a^{\mu} P_{\mu}} \phi(x) e^{-i a^{\mu} P_{\mu}} . \tag{1.144}
\end{equation*}
$$

This, in fact, is simply the transformation property of the quantum field $\phi(x)$ under translation.
Let us now compute the mode expansion for the Heisenberg field. From the mode expansion for the Schrödinger picture operator we find

$$
\begin{equation*}
\phi^{(H)}(t, \vec{x})=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(e^{i H t} a(\vec{p}) e^{-i H t} e^{i \vec{p} \cdot \vec{x}}+e^{i H t} a^{\dagger}(\vec{p}) e^{-i H t} e^{-i \vec{p} \cdot \vec{x}}\right) . \tag{1.145}
\end{equation*}
$$

To simplify this we would like to commute $e^{i H t}$ through $a(\vec{p})$ and $a^{\dagger}(\vec{p})$. Since $[H, a(\vec{p})]=-a(\vec{p}) E_{p}$, we can infer that

$$
\begin{equation*}
H a(\vec{p})=a(\vec{p})\left(H-E_{p}\right) \tag{1.146}
\end{equation*}
$$

and by induction that

$$
\begin{equation*}
H^{n} a(\vec{p})=a(\vec{p})\left(H-E_{p}\right)^{n} . \tag{1.147}
\end{equation*}
$$

Thus

$$
\begin{equation*}
e^{i H t} a(\vec{p})=a(\vec{p}) e^{i\left(H-E_{p}\right) t}, \tag{1.148}
\end{equation*}
$$

which gives the mode expansion in the Heisenberg picture

$$
\begin{equation*}
\phi^{(H)}(t, \vec{x}) \equiv \phi(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(a(\vec{p}) e^{-i p \cdot x}+a^{\dagger}(\vec{p}) e^{i p \cdot x}\right), \tag{1.149}
\end{equation*}
$$

with

$$
\begin{equation*}
p \cdot x=p^{0} x_{0}-\vec{p} \cdot \vec{x}, \quad\left(p^{0}, \vec{p}\right)=\left(E_{p}, \vec{p}\right) . \tag{1.150}
\end{equation*}
$$

In other words, the coefficient of $e^{-i p \cdot x}$ in the mode expansion of the Heisenberg field corresponds to the annihilator and the coefficient of $e^{i p \cdot x}$ to the creator. Note that indeed this mode expansion solves the operator equation of motion (1.142).
For later purposes we also give the inverted expression

$$
\begin{align*}
a(\vec{q}) & =\frac{i}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x e^{i q \cdot x} \stackrel{\leftrightarrow}{\partial_{0}} \phi(x)  \tag{1.151}\\
a^{\dagger}(\vec{q}) & =\frac{-i}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x e^{-i q \cdot x} \stackrel{\leftrightarrow}{\partial_{0}} \phi(x)
\end{align*}
$$

where

- $u(x) \stackrel{\leftrightarrow}{\partial_{0}} v(x):=u(x) \partial_{0} v(x)-\left(\partial_{0} u(x)\right) v(x)$,
- the integrals $\int \mathrm{d}^{3} x$ are evaluated at arbitrary times $t=x^{0}$.

You will check these expressions in the tutorial.

### 1.11 Causality and Propagators

We are finally in a position to come back to the question of causality in a relativistic quantum theory, which in section 1.1 served as one of our two prime motivations to study Quantum Field Theory. We will investigate the problem from two related points of view - via commutators and propagators.

### 1.11.1 Commutators

For causality to hold we need two measurements at spacelike distance not to affect each other. This is guaranteed if any two local observables $O_{1}(x)$ and $O_{2}(y)$ at spacelike separation commute, i.e.

$$
\begin{equation*}
\left[O_{1}(x), O_{2}(y)\right] \stackrel{!}{=} 0 \quad \text { for } \quad(x-y)^{2}<0 \tag{1.152}
\end{equation*}
$$

By a local observable we mean an observable in the sense of quantum mechanics (i.e. a hermitian operator) that is defined locally at a spacetime point $x$, i.e. it depends only on $x$ and at most on a local neighborhood of $x$. Any such local observable is represented by a local (hermitian) operator, by which we mean any local (hermitian) expression of the fundamental operators $\phi(x)$ and $\partial_{\mu} \phi(x)$ such as products or powers series (e.g. exponentials) of the operators. Indeed such operators are local in the sense that they depend only on a local neighborhood of the spacetime point $x$.
To check for (1.152) we thus need to compute the commutator

$$
\begin{equation*}
\Delta(x-y):=[\phi(x), \phi(y)] \tag{1.153}
\end{equation*}
$$

not just at equal times $x^{0}=y^{0}$, but for general times. In Fourier modes we have

$$
\begin{align*}
\Delta(x-y)= & \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}} \int \frac{\mathrm{~d}^{3} q}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{q}}}  \tag{1.154}\\
& \times\left[a(\vec{p}) e^{-i p \cdot x}+a^{\dagger}(\vec{p}) e^{i p \cdot x}, a(\vec{q}) e^{-i q \cdot y}+a^{\dagger}(\vec{q}) e^{i q \cdot y}\right]
\end{align*}
$$

Using the commutation relations for the modes yields

$$
\begin{equation*}
\Delta(x-y)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}\left(e^{-i p \cdot(x-y)}-e^{-i p \cdot(y-x)}\right) . \tag{1.155}
\end{equation*}
$$

Now assume $(x-y)^{2}<0$. Since $\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}$ is Lorentz invariant (see the discussion in section 1.7.2) we can apply a Lorentz transformation such that $\left(x^{0}-y^{0}\right)=0$. Indeed this can be always be achieved if two points are at spacelike distance. This gives

$$
\begin{equation*}
\left.\Delta(x-y)\right|_{(x-y)^{2}<0}=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}\left(e^{i \vec{p} \cdot(\vec{x}-\vec{y})}-e^{-i \vec{p} \cdot(\vec{x}-\vec{y})}\right) \tag{1.156}
\end{equation*}
$$

which is evidently zero because we can change the integration variable from $\vec{p}$ to $-\vec{p}$ in the second term. Consequently, also commutators of the form $\left[\partial_{x^{\mu}} \phi(x), \phi(y)\right]=\partial_{x^{\mu}}[\phi(x), \phi(y)]$ vanish for $x$ and $y$ at spacelike distance. Therefore for all local operators we have

$$
\begin{equation*}
\left[O_{i}(x), O_{j}(y)\right]=0 \quad \text { if }(x-y)^{2}<0 \tag{1.157}
\end{equation*}
$$

This establishes that causality is maintained at the operational level.

## Remark:

In the above we have been able to derive (1.157) because we are working with a free theory, for which the free mode expansion implies $[\phi(x), \phi(y)]=0$ for $(x-y)^{2}<0$. In an interacting theory, such a derivation may not be possible since in this case a free mode expansion is no longer available. More generally, therefore, one must postulate $[\phi(x), \phi(y)]=0$ for $(x-y)^{2}<0$ in form of an axiom of QFT.

## A note on Quantum Entanglement

As we have seen, locality in Quantum Field Theory refers to the fact (or requirement) that local operators commute at spacelike distances. By contrast, the states on which these operators act, i.e. the elements of the Hilbert space, do exhibit non-local behaviour just as in non-relativistic Quantum Mechanics. Indeed the fact that in QFT local operators commute is not at odds with the existence of entangled states, known from Quantum Mechanics. Consider e.g. an entangled 2-spin state $\left|\frac{1}{2},-\frac{1}{2}\right\rangle-$ $\left|-\frac{1}{2}, \frac{1}{2}\right\rangle$. Even though the state is entangled, the spin operator $S_{1}$ for measurement of the spin of particle 1 at $x$ and $S_{2}$ for measurement of the spin of particle 2 at $y$ commute. In particular the expectation value of $S_{2}$ is not changed by measuring $S_{1}$. Therefore by measuring the expectation value of $S_{2}$ we cannot determine if $S_{1}$ was measured and vice versa and the two measurements 'do not affect each other' in this sense. More generally, the states of the Quantum Field Theory are non-local objects. We will understand this better when discussing the Schrödinger Representation of states in the context of path integral quantization in QFT 2. The states are represented as functionals of the field configuration and are necessarily dependent on non-local information.

### 1.11.2 Propagators

Consider now the probability amplitude for a particle emitted at $y$ to propagate to $x$ (where it can be measured). This is given by the propagator

$$
\begin{equation*}
D(x-y):=\langle 0| \phi(x) \phi(y)|0\rangle, \tag{1.158}
\end{equation*}
$$

because $|x\rangle=\phi(x)|0\rangle$ is a position eigenstate. In particular, unlike in the non-relativistic expression (1.3), the time-evolution operator need not be inserted by hand because $|x\rangle$ contains information about the time variable $x^{0}$. By a mode expansion we find

$$
\begin{align*}
D(x-y)= & \iint \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}} \frac{\mathrm{~d}^{3} q}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{q}}}\langle 0|\left(a(\vec{q}) e^{-i q \cdot x}+a^{\dagger}(\vec{q}) e^{i q \cdot x}\right)  \tag{1.159}\\
& \times\left(a(\vec{p}) e^{-i p \cdot y}+a^{\dagger}(\vec{p}) e^{i p \cdot y}\right)|0\rangle .
\end{align*}
$$

Using $a(\vec{p})|0\rangle=0=\langle 0| a^{\dagger}(\vec{q})$ we find that the only contribution is due to the term

$$
\begin{equation*}
e^{-i q \cdot x} e^{i p \cdot x}\langle 0| a(\vec{q}) a^{\dagger}(\vec{p})|0\rangle=e^{-i q \cdot x} e^{i p \cdot x}(2 \pi)^{3} \delta^{(3)}(\vec{q}-\vec{p}) . \tag{1.160}
\end{equation*}
$$

This yields

$$
\begin{equation*}
D(x-y)=\langle 0| \phi(x) \phi(y)|0\rangle=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}} e^{-i p(x-y)} \tag{1.161}
\end{equation*}
$$

We see that $D(x-y)$ is non-zero even for $x$ and $y$ at spacelike distance. On the other hand, in view of (1.156) this cannot affect causality. So how can this be consistent? What saves the day is the observation that

$$
\begin{equation*}
[\phi(x), \phi(y)]=D(x-y)-D(y-x) . \tag{1.162}
\end{equation*}
$$

We can therefore interpret the commutator as describing 2 physical processes, whose quantum probability amplitude apparently cancel each other for $(x-y)^{2}<0$ :

- $D(x-y)$ is the quantum amplitude for a particle to travel from $y \rightarrow x$,
- $D(y-x)$ is the quantum amplitude for a particle to travel from $y \rightarrow x$.

Indeed if $x$ and $y$ are at spacelike distance from each other, there is no Lorentz invariant notion of whether $\left(x^{0}-y^{0}\right)$ is bigger or smaller than zero. Therefore both processes can occur. In the expression for $[\phi(x), \phi(y)]$ these two processes cancel each other in the sense of a destructive quantum mechanical interference.
Even more interestingly we can consider a complex scalar field, which contains the mode operators according to

$$
\begin{align*}
\phi(x) & \sim a(\vec{p}), b^{\dagger}(\vec{p}),  \tag{1.163}\\
\phi^{\dagger}(x) & \sim a^{\dagger}(\vec{p}), b(\vec{p}) .
\end{align*}
$$

For a complex scalar field $[\phi(x), \phi(y)]=0$ for all $x, y$. More interesting is the commutator

$$
\begin{equation*}
\left[\phi(x), \phi^{\dagger}(y)\right]=\langle 0| \phi(x) \phi^{\dagger}(y)|0\rangle-\langle 0| \phi^{\dagger}(y) \phi(x)|0\rangle, \tag{1.164}
\end{equation*}
$$

which vanishes for $(x-y)^{2}<0$ while being in general non-zero otherwise. The first term corresponds to a particle which travels from $y$ to $x$, while the second term corresponds to an anti-particle travelling from $x$ to $y$. Again both processes cancel each other in the expression for the commutator. The important conclusion is that the field formalism saves causality in the QM sense even though the QM probability for a propagation $x \rightarrow y$ itself is non-zero if $(x-y)^{2}<0$. This is why a single-particle approach must fail. The field commutators, on the other hand, know about processes of all possible particles and anti-particles. Thus it is the intrinsic nature of Quantum Field Theory as a multi-particle framework which is responsible for causality.

### 1.11.3 The Feynman-propagator

We will see in the next chapter that an object of crucial importance in QFT is the Feynman-propagator

$$
\begin{equation*}
D_{F}(x-y):=\langle 0| T \phi(x) \phi(y)|0\rangle, \tag{1.165}
\end{equation*}
$$

where

$$
T \phi(x) \phi(y)= \begin{cases}\phi(x) \phi(y) & \text { if } x^{0} \geq y^{0}  \tag{1.166}\\ \phi(y) \phi(x) & \text { if } y^{0}>x^{0}\end{cases}
$$

is the time-ordered product. The Feynman-propagator can be written as

$$
\begin{align*}
D_{F}(x-y)= & \Theta\left(x^{0}-y^{0}\right) \underbrace{\langle 0| \phi(x) \phi(y)|0\rangle}_{D(x-y)}+\Theta\left(y^{0}-x^{0}\right) \underbrace{\langle 0| \phi(y) \phi(x)|0\rangle}_{D(y-x)} \\
= & \left.\Theta\left(x^{0}-y^{0}\right) \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}} e^{-i p \cdot(x-y)}\right|_{p^{0}=+E_{p}}  \tag{1.167}\\
& +\left.\Theta\left(y^{0}-x^{0}\right) \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}} e^{i p \cdot(x-y)}\right|_{p^{0}=+E_{p}}
\end{align*}
$$

Relabeling $\vec{p}$ by $-\vec{p}$ in the second term gives

$$
\begin{equation*}
\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} e^{i \vec{p} \cdot(\vec{x}-\vec{y})}\left[\Theta\left(x^{0}-y^{0}\right) \frac{1}{2 E_{p}} e^{-i E_{p}\left(x^{0}-y^{0}\right)}+\Theta\left(y^{0}-x^{0}\right) \frac{1}{2 E_{p}} e^{i E_{p}\left(x^{0}-y^{0}\right)}\right] . \tag{1.168}
\end{equation*}
$$

The term in brackets can be rewritten as a complex contour integral. As a preparation recall Cauchy's integral formula:

Let $g: U \rightarrow \mathbb{C}$ be a holomorphic function defined in an open subset $U$ of $\mathbb{C}$ and consider a closed disk $D$ with $C=\partial D \subset U$. Then for $z_{0}$ in the interior of $D$ we have

$$
\begin{equation*}
g\left(z_{0}\right)=\frac{1}{2 \pi i} \oint_{C} \frac{g(z)}{z-z_{0}} d z \tag{1.169}
\end{equation*}
$$

In this spirit the first term in $D_{F}(x-y)$ can be written as

$$
\begin{align*}
& \Theta\left(x^{0}-y^{0}\right) \frac{1}{2 E_{p}} e^{-i E_{p}\left(x^{0}-y^{0}\right)} \\
& \quad=-\Theta\left(x^{0}-y^{0}\right) \frac{1}{2 \pi i} \oint_{C_{1}} \mathrm{~d} p^{0} \frac{e^{-i p^{0}\left(x^{0}-y^{0}\right)}}{\left(p^{0}-E_{p}\right)\left(p^{0}+E_{p}\right)}, \tag{1.170}
\end{align*}
$$

where

- the poles are avoided in $\epsilon$-surroundings as drawn in the picture.
- we close the contour in the lower half-plane such as to pick up the residue at $p^{0}=+E_{p}$ and
- the integral is clockwise, which explains the overall minus sign.


Figure 1.1: Contour $\mathcal{C}_{1}$.

It is important to note that this result holds for any contour that encloses the poles as drawn. In particular, since $x^{0}-y^{0}>0$ the integral along the lower half-plane asymptotically vanishes if we choose to deform the contour to infinity, corresponding to $R \rightarrow \infty$. For the second term in (1.168) we can similarly write

$$
\begin{align*}
& \Theta\left(y^{0}-x^{0}\right) \frac{1}{2 E_{p}} e^{i E_{p}\left(x^{0}-y^{0}\right)} \\
& \quad=-\Theta\left(y^{0}-x^{0}\right) \frac{1}{2 \pi i} \oint_{C_{2}} \mathrm{~d} p^{0} \frac{e^{-i p^{0}\left(x^{0}-y^{0}\right)}}{\left(p^{0}-E_{p}\right)\left(p^{0}+E_{p}\right)} . \tag{1.171}
\end{align*}
$$

This time $\mathcal{C}_{2}$ runs counter-clockwise, but picks up the pole at $p_{0}=-E_{p}$, which again yields an overall minus.
Now, as stressed above both expressions hold for any $R>E_{p}$, but if $R \rightarrow \infty$, then the integral in the lower-/upper halfplane each vanishes due to the appearance of $\Theta\left(x^{0}-y^{0}\right)$ and $\Theta\left(y^{0}-x^{0}\right)$, respectively. It is at this place that the time ordering becomes crucial.
We can therefore evaluate both integrals for $R \rightarrow \infty$, add them with the help of $1=\Theta\left(x^{0}-y^{0}\right)+$ $\Theta\left(y^{0}-x^{0}\right)$ and arrive at

$$
\begin{equation*}
D_{F}(x-y)=\oint_{C} \frac{\mathrm{~d}^{4} p}{(2 \pi)^{4}} \frac{i}{p^{2}-m^{2}} e^{-i p \cdot(x-y)}, \tag{1.172}
\end{equation*}
$$



Figure 1.2: Contour $\mathcal{C}_{2}$.


Figure 1.3: Contour $C$ of the Feynman propagator.
where $\left(p^{0}-E_{p}\right)\left(p^{0}+E_{p}\right)=p^{2}-m^{2}$ was used. In this expression the contour integral in $p^{0}$ must be taken along the path $C$ shown above.
Since this is a complex contour integral, all that matters is the relative position of the contour to the poles. Therefore we can equivalently pick the contour $C$ on top of the real axis but shift the poles e.g. by an amount $\pm i \tilde{\epsilon} / E_{p}$ in the limit $\tilde{\epsilon} \rightarrow 0$.


Figure 1.4: Shifted poles.

This modifies the denominators as

$$
\begin{equation*}
p^{0}=-E_{p}+i \tilde{\epsilon} / E_{p} \quad \text { and } \quad p^{0}=E_{p}-i \tilde{\epsilon}_{p} / E_{p} . \tag{1.173}
\end{equation*}
$$

This must be combined with taking the limit $\tilde{\epsilon} \rightarrow 0$ after performing the integral. With this understood and using furthermore $\left.\left(p^{0}-\left(E_{p}-i \tilde{\epsilon} / E_{p}\right)\right)\left(p^{0}+\overline{\left(E_{p}\right.}-i \tilde{\epsilon} / E_{p}\right)\right)=\left(p^{0}\right)^{2}-E_{p}^{2}+2 i \tilde{\epsilon}+\epsilon^{2} / E_{p}^{2}=$ $p^{2}-m^{2}+i \epsilon$ with $\epsilon=2 i \tilde{\epsilon}-i \tilde{\epsilon} / E_{p}$, the Feynman propagator can be written as the integral

$$
\begin{equation*}
D_{F}(x-y)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i}{p^{2}-m^{2}+i \epsilon} e^{-i p \cdot(x-y)} \tag{1.174}
\end{equation*}
$$

with the $p^{0}$ integration along the real axis and with the limit $\epsilon \rightarrow 0$ after performing the integral.

### 1.11.4 Propagators as Green's functions

Direct computation reveals that $D_{F}(x-y)$ is a Green's function for the Klein-Gordon equation

$$
\begin{equation*}
\left(\partial_{x}^{2}+m^{2}\right) D_{F}(x-y)=-i \delta^{(4)}(x-y) \tag{1.175}
\end{equation*}
$$

The general solution to the equation

$$
\begin{equation*}
\left(\partial^{2}+m^{2}\right) \Delta(x)=-i \delta^{(4)}(x) \tag{1.176}
\end{equation*}
$$

is found by Fourier transforming both sides as

$$
\begin{equation*}
\Delta(x)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \tilde{\Delta}(p) e^{-i p \cdot x}, \delta^{(4)}(x)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} e^{-i p \cdot x} \tag{1.177}
\end{equation*}
$$

and noting that the Klein-Gordon equation becomes an algebraic equation for the Fourier transforms,

$$
\begin{equation*}
\left(-p^{2}+m^{2}\right) \tilde{\Delta}(p)=-i \Rightarrow \tilde{\Delta}(p)=\frac{i}{p^{2}-m^{2}} \tag{1.178}
\end{equation*}
$$

Special care must now be applied in performing the Fourier backtransformation for this solution due to the appearance of the two poles at $p^{0}= \pm \sqrt{E_{\vec{p}}^{2}}$. Any consistent prescription to avoid divergences in performing the contour integral in $p^{0}$ leads to a solution to the original equation (1.175).
In fact, there are $2 \times 2$ different ways to evaluate the contour integral in $p^{0}$ by avoiding the two poles in the upper and lower half-plane. As will be shown in the tutorials, if we avoid both poles along a contour in the upper half-plane, the solution is the retarded Green's function

$$
\begin{align*}
D_{R}(x-y) & =\Theta\left(x^{0}-y^{0}\right)[D(x-y)-D(y-x)]  \tag{1.179}\\
& \equiv \Theta\left(x^{0}-y^{0}\right)\langle 0|[\phi(x), \phi(y)]|0\rangle
\end{align*}
$$

while avoiding both poles in the lower half-plane yields the the advanced Green's function

$$
\begin{equation*}
D_{A}(x-y)=\Theta\left(y^{0}-x^{0}\right)[D(x-y)-D(y-x)] . \tag{1.180}
\end{equation*}
$$



Figure 1.5: The retarded contour.


Figure 1.6: The advanced contour.

Note that $D_{R}$ and $D_{A}$ appear also in classical field theory in the context of constructing solutions to the inhomogenous Klein-Gordon equation, where they propagate the inhomogeneity forward ( $D_{R}$ ) and backward $\left(D_{A}\right)$ in time. In particular the classical version of causality is the statement that $D_{R}$ and $D_{A}$ vanish if $(x-y)^{2}<0$, which we proved from a different perspective before.
By contrast $D_{F}(x-y)$ is the solution corresponding to the contour described previously in this section. It does not appear in classical field theory. The reason is that $D_{F}(x-y)$ propagates 'positive frequency modes' $e^{-i p x}$ forward and 'negative frequency modes' $e^{i p x}$ backward in time (see (1.168)). Remember that $D_{F}(x-y)$ is non-vanishing, even for $x$ and $y$ at spacelike distance. Finally there is a fourth prescription which has no particularly important interpretation in classical or quantum field theory.

## Chapter 2

## Interacting scalar theory

### 2.1 Introduction

So far we have considered a free scalar theory

$$
\begin{equation*}
\mathcal{L}_{0}=\frac{1}{2}(\partial \phi)^{2}-V_{0}(\phi) \tag{2.1}
\end{equation*}
$$

with $V_{0}(\phi)=\frac{1}{2} m_{0}^{2} \phi^{2}$.

- The theory is exactly solvable: The Hilbert space is completely known as the Fock space of multi-particle states created from the vacuum $|0\rangle$.
- There are no interactions between the particles.

Interactions are described in QFT by potentials $V(\phi)$ beyond quadratic order. We can think of $V(\phi)$ as a formal power series in $\phi$,

$$
\begin{equation*}
V(\phi)=\underbrace{\frac{1}{2} m_{0}^{2} \phi^{2}}_{V_{0}}+\underbrace{\frac{1}{3!} g \phi^{3}+\frac{1}{4!} \lambda \phi^{4}+\ldots,}_{V_{\mathrm{int}}} \tag{2.2}
\end{equation*}
$$

and decompose the full Langrangian as

$$
\begin{equation*}
\mathcal{L}=\mathcal{L}_{0}+\mathcal{L}_{\mathrm{int}}, \quad \mathcal{L}_{\mathrm{int}}=-V_{\mathrm{int}} \tag{2.3}
\end{equation*}
$$

This yields the decomposition of the full Hamiltonian as $H=H_{0}+H_{\text {int }}$. Introducing interaction terms leads to a number of important changes in the theory:

- The Hilbert space is different from the Hamiltonian of the free theory.
- This is true already for the vacuum, i.e. the full Hamiltonian $H$ has a ground state $|\Omega\rangle$ different from the ground state $|0\rangle$ of the free Hamiltonian $H_{0}$ :

$$
\begin{align*}
& |0\rangle \leftrightarrow \text { vacuum of } H_{0}: H_{0}|0\rangle=E_{0}|0\rangle  \tag{2.4}\\
& |\Omega\rangle \leftrightarrow \text { vacuum of } H: H|\Omega\rangle=E_{\Omega}|\Omega\rangle
\end{align*}
$$

with $|\Omega\rangle \neq|0\rangle$.

- The mass of the momentum eigenstates of $H$ does no longer equal the parameter $m_{0}$ that appears in $\mathcal{L}_{0}$.
- Bound states may exist in the spectrum.
- The states interact.

The dream of Quantum Field Theorists is to find the exact solution of a non-free QFT, i.e. find the exact spectrum and compute all interactions exactly. This has only been possible so far for very special theories with a lot of symmetry, e.g. certain 2-dimensional QFTs with conformal invariance (Conformal Field Theory), or certain 4-dimensional QFTs with enough supersymmetry.
However, for small coupling parameters such as $g$ and $\lambda$ in $V(\phi)$ we can view the higher terms as small perturbations and apply perturbation theory. Note that depending on the mass dimension of the couplings it must be specified in what sense these parameters must be small, but e.g. for the dimensionless parameter $\lambda$ this would mean that $\lambda \ll 1$ for perturbation theory to be applicable.
Before dealing with interactions in such a perturbative approach, however, we will be able to establish a number of non-trivial important results on the structure of the spectrum and interactions in a nonperturbative fashion.

### 2.2 Källén-Lehmann spectral representation

We take a first look at the spectrum of an interacting real scalar field theory in a manner valid for all types of interactions and without relying on perturbation theory. As a consequence of Lorentz invariance the Hamiltonian and the 3-momentum operator must of course still commute, $[H, \vec{P}]=0^{1}$, and can thus be diagonalised simultaneously. By $\left|\lambda_{\vec{p}}\right\rangle$ we denote such an eigenstate of $H$ and $\vec{P}$ in the full theory such that

$$
\begin{align*}
H\left|\lambda_{\vec{p}}\right\rangle & =E_{p}(\lambda)\left|\lambda_{\vec{p}}\right\rangle, \\
\vec{P}\left|\lambda_{\vec{p}}\right\rangle & =\vec{p}\left|\lambda_{\vec{p}}\right\rangle . \tag{2.5}
\end{align*}
$$

Each $\left|\lambda_{\vec{p}}\right\rangle$ is related via a Lorentz boost with the corresponding state at rest, called $\left|\lambda_{0}\right\rangle$. We can have the following types of $\left|\lambda_{\vec{p}}\right\rangle$ :

- 1-particle states $\left|1_{\vec{p}}\right\rangle$ with $E_{p}=\sqrt{\vec{p}^{2}+m^{2}}$ and rest mass $m$. Remember that $m$ is no more identical to $m_{0}$ in $\mathcal{L}_{0}$.
- Bound states with no analogue in the free theory.
- 2- and $N$-particle states formed out of 1-particle and the bound states. In this case we take $\vec{p}$ to be the centre-of-mass momentum of the multi-particle state.

[^5]All these states are created from the vacuum $|\Omega\rangle$. The crucial difference to the free theory is, though, that $\phi(x)$ cannot simply be written as a superposition of its Fourier amplitudes $a(\vec{p})$ and $a^{\dagger}(\vec{p})$ because it does not obey the free equation of motion, i.e.

$$
\begin{equation*}
\left(\partial^{2}+m^{2}\right) \phi \neq 0 . \tag{2.6}
\end{equation*}
$$

Rather,

$$
\begin{equation*}
\left(\partial^{2}+m^{2}\right) \phi=j \tag{2.7}
\end{equation*}
$$

for a suitable current $j$. Thus, acting with $\phi$ on $|\Omega\rangle$ does not simply create a 1-particle state as in the free theory.
We will make frequent use of the completeness relation of this Hilbert space,

$$
\begin{equation*}
\mathbb{1}=|\Omega\rangle\langle\Omega|+\sum_{\lambda} \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}(\lambda)}\left|\lambda_{\vec{p}}\right\rangle\left\langle\lambda_{\vec{p}}\right| . \tag{2.8}
\end{equation*}
$$

Here the formal sum over $\lambda$ includes the sum over the 1 -particle state, over all types of bound states as well as over all multi-particle states, while the integral $\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{E_{p}(\lambda)}$ refers to the centre-of-mass momentum of a state of species $\lambda$. In particular, since specifying a multi-particle state requires specifying the relative momenta of the individual states (in addition to the centre-of-mass momentum $\vec{p}$ ), the sum over $\lambda$ is really a sum over a continuum of states.
Our first goal is to compute the interacting Feynman-propagator

$$
\begin{equation*}
\langle\Omega| T \phi(x) \phi(y)|\Omega\rangle \tag{2.9}
\end{equation*}
$$

and to establish a physical interpretation for it. Even though we cannot rely on the mode expansion of the field any more, we will be able to make a great deal of progress with the help of two tricks.

- First we insert $\mathbb{1}$ between $\phi(x)$ and $\phi(y)^{2}$. We first ignore time ordering. Then, without loss of generality we can assume that

$$
\begin{equation*}
\langle\Omega| \phi(x)|\Omega\rangle=0, \tag{2.10}
\end{equation*}
$$

because since

$$
\begin{equation*}
\phi(x)=e^{i x^{\mu} P_{\mu}} \phi(0) e^{-i x^{\mu} P_{\mu}} \tag{2.11}
\end{equation*}
$$

(cf. (1.144)) and $P^{\mu}|\Omega\rangle=0$ we have

$$
\begin{equation*}
\langle\Omega| \phi(x)|\Omega\rangle=\langle\Omega| \phi(0)|\Omega\rangle \forall x . \tag{2.12}
\end{equation*}
$$

So if $c \equiv\langle\Omega| \phi(0)|\Omega\rangle \neq 0$ we simply redefine $\phi \rightarrow \phi-c$ to achieve (2.10). Therefore

$$
\begin{equation*}
\langle\Omega| \phi(x) \mathbb{1} \phi(y)|\Omega\rangle=\sum_{\lambda} \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}(\lambda)}\langle\Omega| \phi(x)\left|\lambda_{\vec{p}}\right\rangle\left\langle\lambda_{\vec{p}}\right| \phi(y)|\Omega\rangle . \tag{2.13}
\end{equation*}
$$

[^6]Now

$$
\begin{equation*}
\langle\Omega| \phi(x)\left|\lambda_{\vec{p}}\right\rangle=\langle\Omega| e^{i P \cdot x} \phi(0) e^{-i P \cdot x}\left|\lambda_{\vec{p}}\right\rangle, \tag{2.14}
\end{equation*}
$$

and with $\langle\Omega| e^{i P \cdot x}=\langle\Omega|$ and $e^{-i P \cdot x}\left|\lambda_{\vec{p}}\right\rangle=\left|\lambda_{\vec{p}}\right\rangle e^{-i p \cdot x}$ we find

$$
\begin{equation*}
\langle\Omega| \phi(x)\left|\lambda_{\vec{p}}\right\rangle=\langle\Omega| \phi(0)\left|\lambda_{\vec{p}}\right\rangle e^{-i p \cdot x} . \tag{2.15}
\end{equation*}
$$

- The next trick is to relate $\left|\lambda_{\vec{p}}\right\rangle$ to $\left|\lambda_{0}\right\rangle$ by a Lorentz boost. To this end we investigate the transformation behaviour of a scalar field under a Lorentz transformation

$$
\begin{equation*}
x \mapsto x^{\prime}=\Lambda x . \tag{2.16}
\end{equation*}
$$

In the spirit of Quantum Mechanics the action of the Lorentz group is represented on the Hilbert space in terms of a unitary operator $U(\Lambda)$ such that all states transform like

$$
\begin{equation*}
|\alpha\rangle \mapsto\left|\alpha^{\prime}\right\rangle=U(\Lambda)|\alpha\rangle . \tag{2.17}
\end{equation*}
$$

What is new to us is that also the transformation of the field is determined in terms of $U(\Lambda)$. More precisely, the scalar field transforms as

$$
\begin{equation*}
\phi(x) \mapsto \phi^{\prime}\left(x^{\prime}\right)=\phi\left(x\left(x^{\prime}\right)\right) \tag{2.18}
\end{equation*}
$$

with $\phi^{\prime}\left(x^{\prime}\right)=U^{-1}(\Lambda) \phi\left(x^{\prime}\right) U(\Lambda)$, i.e.

$$
\begin{equation*}
U^{-1}(\Lambda) \phi\left(x^{\prime}\right) U(\Lambda)=\phi(x) . \tag{2.19}
\end{equation*}
$$

To see this we start with the familiar transformation of a classical scalar field under a Lorentz transformation given by

$$
\begin{equation*}
\phi \mapsto \phi^{\prime}\left(x^{\prime}\right)=\phi(x) . \tag{2.20}
\end{equation*}
$$

We now need to find the analogue of this equation for operator-valued fields. The analogue of the classical value of $\phi(x)$ is the matrix element $\langle\alpha| \phi(x)|\beta\rangle$ evaluated on a basis of the Hilbert space. The transformed field $\phi^{\prime}(x)$ then corresponds to transformed matrix elements $\left\langle\alpha^{\prime}\right| \phi(x)\left|\beta^{\prime}\right\rangle$. Thus the classical relation (2.20) translates into

$$
\begin{equation*}
\underbrace{\left\langle\alpha^{\prime}\right| \phi\left(x^{\prime}\right)\left|\beta^{\prime}\right\rangle}_{\langle\alpha| U^{-1} \phi\left(x^{\prime}\right) U|\beta\rangle}=\langle\alpha| \phi(x)|\beta\rangle, \tag{2.21}
\end{equation*}
$$

for all states $|\alpha\rangle$ and $|\beta\rangle$.


$$
\begin{equation*}
\left|\lambda_{\vec{p}}\right\rangle=U^{-1}\left|\lambda_{0}\right\rangle . \tag{2.22}
\end{equation*}
$$

We can then further manipulate (2.15) by writing

$$
\begin{equation*}
\langle\Omega| \phi(x)\left|\lambda_{\vec{p}}\right\rangle=\underbrace{\langle\Omega| U^{-1}}_{\langle\Omega|} \underbrace{U \phi(0) U^{-1}}_{\phi(0)} \underbrace{U\left|\lambda_{\vec{p}}\right\rangle}_{\left|\lambda_{0}\right\rangle} e^{-i p \cdot x} . \tag{2.23}
\end{equation*}
$$



Figure 2.1: Spectral function.

Thus the Feynman propagator without time ordering can be expressed as

$$
\begin{equation*}
\left.\langle\Omega| \phi(x) \phi(y)|\Omega\rangle=\sum_{\lambda} \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}(\lambda)} e^{-i p \cdot(x-y)}|\langle\Omega| \phi(0)| \lambda_{0}\right\rangle\left.\right|^{2}, \tag{2.24}
\end{equation*}
$$

which is to be compared with the free scalar result (1.161). Including time ordering, we can peform the same manipulations for the integral as in the free theory and therefore conclude

$$
\begin{equation*}
\left.\langle\Omega| T \phi(x) \phi(y)|\Omega\rangle=\sum_{\lambda} \int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i}{p^{2}-m_{\lambda}^{2}+i \epsilon} e^{-i p \cdot(x-y)}|\langle\Omega| \phi(0)| \lambda_{0}\right\rangle\left.\right|^{2} \tag{2.25}
\end{equation*}
$$

One can define

$$
\begin{equation*}
D_{F}\left(x-y ; M^{2}\right):=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i}{p^{2}-M^{2}+i \epsilon} e^{-i p \cdot(x-y)} \tag{2.26}
\end{equation*}
$$

to write

$$
\begin{equation*}
\langle\Omega| T \phi(x) \phi(y)|\Omega\rangle=\int_{0}^{\infty} \frac{\mathrm{d} M^{2}}{2 \pi} \rho\left(M^{2}\right) D_{F}\left(x-y ; M^{2}\right) \tag{2.27}
\end{equation*}
$$

in terms of the spectral function

$$
\begin{equation*}
\left.\rho\left(M^{2}\right)=\sum_{\lambda} 2 \pi \delta\left(M^{2}-m_{\lambda}^{2}\right)|\langle\Omega| \phi(0)| \lambda_{0}\right\rangle\left.\right|^{2}, \tag{2.28}
\end{equation*}
$$

which has a typical form like in Figure 2.1. It is crucial to appreciate that the 1-particle states leads to an isolated $\delta$-function peak around $M^{2}=m^{2}$. Therefore below $M^{2} \cong(2 m)^{2}$ or $M^{2} \cong m_{\text {bound }}^{2}$ the spectral function takes the form

$$
\begin{equation*}
\rho\left(M^{2}\right)=2 \pi \delta\left(M^{2}-m^{2}\right) Z . \tag{2.29}
\end{equation*}
$$

Here we have defined the field-strength or wavefunction renormalisation

$$
\begin{equation*}
\left.Z=|\langle\Omega| \phi(0)| 1_{0}\right\rangle\left.\right|^{2} \tag{2.30}
\end{equation*}
$$

where $\left|1_{0}\right\rangle$ is a 1-particle state at rest.

Consider now the Fourier-transformation

$$
\begin{align*}
\int \mathrm{d}^{4} x e^{i p \cdot(x-y)}\langle\Omega| T \phi(x) \phi(y)|\Omega\rangle & =\int_{0}^{\infty} \frac{\mathrm{d} M^{2}}{2 \pi} \rho\left(M^{2}\right) \frac{i}{p^{2}-M^{2}+i \epsilon}  \tag{2.31}\\
& =\frac{i Z}{p^{2}-m^{2}+i \epsilon}+\int_{m_{\text {bound }}^{2}}^{\infty} \frac{\mathrm{d} M^{2}}{2 \pi} \rho\left(M^{2}\right) \frac{i}{p^{2}-M^{2}+i \epsilon} .
\end{align*}
$$

This identifies the 1-particle state as the first analytic pole at $m^{2}$. This is an important observation:
The mass-square $m^{2}$ of the particle is the location of the lowest-lying pole of the Fourier transformed propagator.
I.e. computation of the propagator gives us, among other things, a way to read off the mass of the particle. Note furthermore that
in the complex $\mathrm{p}^{\wedge} 2$ plane?

- bound states appear at higher isolated poles and
- $N$-particle states give rise to a branch cut beginning at $p^{2}=4 m^{2}$.

The field-strength renormalisation $Z$ was 1 in free theory because $\phi(0)$ just creates the free particle from vacuum. In an interacting theory

$$
\begin{equation*}
\left.|\langle\Omega| \phi(0)| 1_{0}\right\rangle \mid=\sqrt{Z}<1 \tag{2.32}
\end{equation*}
$$

because $\phi$ creates not only 1-particle states and thus the overlap with the 1-particle states is smaller. In fact, by exploiting the properties of the spectral function one can prove formally that
$Z=1$ if and only if the theory is free.
We give a guided tour through this proof in the tutorials.

### 2.3 S-matrix and asymptotic in/out-states

We now consider scattering of incoming states $|i\rangle$ to outgoing states $|f\rangle$ with the aim of computing the QM transition amplitude, i.e. the probability amplitude for scattering of $|i\rangle$ to $|f\rangle$. The process is formulated in terms of the theory of asymptotic in- and out-states.

- In the asymptotic past, $t \rightarrow-\infty$, the in-states $\mid i$, in $\rangle$ are described as distinct wave-packets corresponding to well-separated single particle states. Being far apart for $t \rightarrow-\infty$, they travel freely as individual states. This is a consequence of locality of the interactions, which we assume in the sequel.
- As these states approach each other, they start to "feel each other", interact and scatter into the final states $|f\rangle$.
- For $t \rightarrow \infty$ these final states are again asymptotically free and well-separated 1-particle states. The concept of free asymptotic in/out states is formalised by the so-called in- and out-fields $\phi_{\text {in }}$ and $\phi_{\text {out }}$ with the following properties:
- The in-state $\mid i$, in $\rangle$ is created from the asymptotic vaccum $\mid$ vac, in $\rangle$ by action of $\phi_{\text {in }}$ as $t \rightarrow-\infty$. We will see that

the vacuum of the interacting theory.
not logical. this is being invented.
- $\mid i$, in $\rangle$ has $E=\sqrt{p^{2}+m^{2}}$ with $m$ the value of the 1-particle pele in the Feynman propagator of the full interacting theory. In particular $m \neq m_{0}$. Therefore, $\phi_{\text {in }}$ is a free field obeying the free Klein-Gordon-equation, but with the full mass $m \neq m_{0}$,

$$
\begin{equation*}
\left(\partial^{2}+m^{2}\right) \phi_{\text {in }}=0 . \tag{2.34}
\end{equation*}
$$

It is thus possible to expand $\phi_{\text {in }}$ in terms of $a_{\text {in }}(\vec{p})$ and $a_{\text {in }}^{\dagger}(\vec{p})$ so that

$$
\begin{equation*}
\phi_{\mathrm{in}}(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(a_{\mathrm{in}}(\vec{p}) e^{-i p \cdot x}+a_{\mathrm{in}}^{\dagger}(\vec{p}) e^{i p \cdot x}\right), \tag{2.35}
\end{equation*}
$$

where $p^{0}=\sqrt{\vec{p}^{2}+m^{2}}$.

- $\phi_{\text {in }}$ satisfies the following relation to the interacting field $\phi$ : Asymptotically for $t \rightarrow-\infty$ the above logic suggests identifying $\phi_{\text {in }}(x)$ with $\phi(x)$, at least in weak sense that their their matrix elements with a basis of the Hilbert space must agree in a suitable manner. We therefore make the ansatz

$$
\begin{equation*}
\phi \rightarrow C \phi_{\mathrm{in}} \tag{2.36}
\end{equation*}
$$

in the weak sense only, i.e.

$$
\begin{equation*}
\langle\alpha| \phi|\beta\rangle \rightarrow C\langle\alpha| \phi_{\text {in }}|\beta\rangle \tag{2.37}
\end{equation*}
$$

for all $|\alpha\rangle$ and $|\beta\rangle$ as $t \rightarrow-\infty$. With this input one can show (see Examples Sheet 5 for the proof) that

$$
\begin{equation*}
\left\langle 1_{\vec{p}}\right| \phi(0)|\Omega\rangle=C\left\langle 1_{\vec{p}}\right| \phi_{\text {in }}(0)|\Omega\rangle . \tag{2.38}
\end{equation*}
$$

Since by construction $\left\langle 1_{\vec{p}}\right| \phi_{\text {in }}(0)|\Omega\rangle=1$ this identifies $C=\sqrt{Z}$ with $Z$ the wavefunction renormalisation of the full theory as defined in (2.30). Thus

$$
\begin{equation*}
\langle\alpha| \phi|\beta\rangle \rightarrow \sqrt{Z}\langle\alpha| \phi_{\text {in }}|\beta\rangle \tag{2.39}
\end{equation*}
$$

as $t \rightarrow-\infty$. Note that this is really true only in this weak sense. What does not hold in an interacting field theory is that operator products (i.e. powers of $\phi(x)$ ) approach corresponding products of $\phi_{\text {in }}(x) .^{3}$

[^7]- $\phi_{\text {out }}$ has of course analogous properties as $t \rightarrow+\infty$.

Our considerations can be summarised as follows: $\phi_{\text {in/out }}$ are free fields with single particle states of energy $E_{p}=\sqrt{p^{2}+m^{2}}$, with the mass $m$ of the full theory. We can think of switching off all interactions of the theory as $t \rightarrow \mp \infty$ except for self-interactions of the field. This leads to mass $m \neq m_{0}$ and $Z \neq 1$. We will be able to understand what is meant by these self-interactions very soon when discussing the resummed propagator, and it will become clear then that indeed it is imperative to take into account $m \neq m_{0}$ and $Z \neq 1$ for $\phi_{\text {in/out }}$.

The Hilbert spaces of asymptotic in- and out-states are isomorphic Fock spaces. Thus there exists an operator $S$ which maps the out-states onto the in-states, i.e.


On Examples Sheet 5 we prove the following properties of $S$ :

- $S$ is unitary, i.e $S^{\dagger}=S^{-1}$,
- $\phi_{\text {in }}(x)=S \phi_{\text {out }}(x) S^{-1}$,
- $\mid$ vac, in $\rangle=\mid$ vac,out $\rangle=|\Omega\rangle$ and $S|\Omega\rangle=|\Omega\rangle$.

Our aim is to compute the transition amplitude

$$
\begin{equation*}
\langle f, \text { out }| i, \text { in }\rangle=\underbrace{\langle f, \text { in }| S \mid i, \text { in }\rangle}_{S \text {-matrix element }}, \tag{2.41}
\end{equation*}
$$

so that $\mid\langle f$, in $| S \mid i$, in $\rangle\left.\right|^{2}$ it the probability for scattering from the initial states to the final states.

### 2.4 The LSZ reduction formula

Let us now compute the $S$-matrix element

$$
\begin{equation*}
\left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| q_{1}, \ldots, q_{r}, \text { in }\right\rangle \tag{2.42}
\end{equation*}
$$

for scattering of asymptotic in- and out-states of definite momenta. One can think of this as the building block to describe scattering of asymptotically localised wave-packets $\left|f_{\text {in }}\right\rangle=\int \mathrm{d}^{3} p f(\vec{p})\left|p_{\text {in }}\right\rangle$.

- We first use the definition of $\mid q_{i}$, in $\rangle$ in terms of the creation operator of the in-field

$$
\begin{equation*}
\phi_{\mathrm{in}}(x)=\int \frac{\mathrm{d}^{3} k}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{k}}}\left(a_{\mathrm{in}}(\vec{k}) e^{-i k \cdot x}+a_{\mathrm{in}}^{\dagger}(\vec{k}) e^{i k \cdot x}\right) \tag{2.43}
\end{equation*}
$$

given by

$$
\begin{equation*}
\left|q_{i}, \mathrm{in}\right\rangle=\sqrt{2 E_{q_{i}}} a_{\mathrm{in}}^{\dagger}\left(q_{i}\right)|\Omega\rangle . \tag{2.44}
\end{equation*}
$$

We will use that

$$
\begin{align*}
& a_{\text {in }}(\vec{q})=\left.\frac{i}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x e^{i q \cdot x} \stackrel{\leftrightarrow}{\partial_{0}} \phi_{\text {in }}(x)\right|_{x^{0}=t}  \tag{2.45}\\
& a_{\mathrm{in}}^{\dagger}(\vec{q})=\left.\frac{-i}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x e^{-i q \cdot x} \stackrel{\leftrightarrow}{\partial_{0}} \phi_{\text {in }}(x)\right|_{x^{0}=t}
\end{align*}
$$

(cf. (1.151)), where the integral can be evaluated at arbitrary time $t$. This allows us to trade $\left|q_{i}\right\rangle$ by $\phi_{\text {in }}(x)$ as follows:

$$
\begin{array}{r}
\left.\left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| q_{1}, \ldots, q_{r}, \text { in }\right\rangle=\sqrt{2 E_{q_{1}}}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| a_{\mathrm{in}}^{\dagger}\left(\overrightarrow{q_{1}}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle \\
\left.\left.=\frac{1}{i} \int \mathrm{~d}^{3} x e^{-i q_{1} \cdot x} \stackrel{\leftrightarrow}{\partial_{0}}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi_{\text {in }}(t, \vec{x}) \right\rvert\, q_{2}, \ldots, q_{r}, \text { in }\right\rangle\left.\right|_{x^{0}=t} \tag{2.46}
\end{array}
$$

- Since $t$ is arbitrary, we can take $t \rightarrow-\infty$ because in that limit we can make use of the relation

$$
\begin{equation*}
\lim _{t \rightarrow-\infty}\left\langle 1_{\vec{p}}\right| \phi_{\mathrm{in}}(t, \vec{x})|\Omega\rangle=\lim _{t \rightarrow-\infty} Z^{-1 / 2}\left\langle 1_{\vec{p}}\right| \phi(t, \vec{x})|\Omega\rangle, \tag{2.47}
\end{equation*}
$$

or more generally

$$
\begin{equation*}
\left.\left.\lim _{t \rightarrow-\infty}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi_{\text {in }}(t, \vec{x}) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle=\lim _{t \rightarrow-\infty} Z^{-1 / 2}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi(x) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle \tag{2.48}
\end{equation*}
$$

This leads to

$$
\begin{align*}
& \left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| q_{1}, \ldots, q_{r}, \text { in }\right\rangle \\
& =\lim _{t \rightarrow-\infty} Z^{-1 / 2} \underbrace{\left.\left.\frac{1}{i} \int \mathrm{~d}^{3} x e^{-i q_{1} \cdot x} \stackrel{\leftrightarrow}{\partial_{0}}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi(x) \right\rvert\, q_{2}, \ldots, q_{r}, \text { in }\right\rangle}_{\equiv \int \mathrm{d}^{3} x f(t, \vec{x})} \tag{2.49}
\end{align*}
$$

- Our next aim is to let $\phi(x)$ act from the right on the out-states in order to annihilate one of the states. To this end we need to relate the above matrix element with $t \rightarrow-\infty$ to a matrix element in the limit $t \rightarrow \infty$, where we can re-express $\phi(x)$ by $\phi_{\text {out }}(x)$. We can do so by exploiting that for all functions $f(t, \vec{x})$

$$
\begin{equation*}
\left(\lim _{t \rightarrow \infty}-\lim _{t \rightarrow-\infty}\right) \int \mathrm{d}^{3} x f(t, \vec{x})=\underbrace{\lim _{\substack{t_{f} \rightarrow \infty \\ t_{t \rightarrow-\infty}}}^{\int_{t_{i}}} \mathrm{~d} t \frac{\partial}{\partial t} \int \mathrm{~d}^{3} x f(t, \vec{x})}_{\equiv \int \mathrm{d}^{4} x \partial_{0} f(x)} \tag{2.50}
\end{equation*}
$$

i.e.

$$
\begin{equation*}
\lim _{t \rightarrow-\infty} \int \mathrm{d}^{3} x f(t, \vec{x})=\lim _{t \rightarrow+\infty} \int \mathrm{d}^{3} x f(t, \vec{x})-\int \mathrm{d}^{4} x \partial_{0} f(x) \tag{2.51}
\end{equation*}
$$

Therefore we can write the overlap of the in- and out-states as

$$
\begin{equation*}
\left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| q_{1}, \ldots, q_{r}, \text { in }\right\rangle=A-B \tag{2.52}
\end{equation*}
$$

with

$$
\begin{equation*}
\left.B=\int \mathrm{d}^{4} x Z^{-1 / 2} \partial_{0}\left[\left.e^{-i q_{1} \cdot x} \frac{1}{i} \stackrel{\leftrightarrow}{\partial_{0}}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi(x) \right\rvert\, q_{2}, \ldots, q_{r}, \text { in }\right\rangle\right] \tag{2.53}
\end{equation*}
$$

and

$$
\begin{align*}
A & =\lim _{t \rightarrow \infty} \int \mathrm{~d}^{3} x \frac{1}{i} e^{-i q_{1} \cdot x} \stackrel{\leftrightarrow}{\partial_{0}} \underbrace{\left.Z^{-1 / 2}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi(x) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle}_{\text {for } \left.t \rightarrow \infty:\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi_{\text {out }}(x) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle} \\
& \left.=\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| a_{\text {out }}^{\dagger}\left(\vec{q}_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle \sqrt{2 E_{q_{1}}} . \tag{2.54}
\end{align*}
$$

Altogether

$$
\begin{align*}
\left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| q_{1}, \ldots, q_{r}, \text { in }\right\rangle= & \left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| a_{\mathrm{out}}^{\dagger}\left(\vec{q}_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle \sqrt{2 E_{q_{1}}}  \tag{2.55}\\
& \left.+i \int \mathrm{~d}^{4} x Z^{-1 / 2} \partial_{0}\left[e^{-i q_{1} \cdot x} \stackrel{\leftrightarrow}{\partial_{0}}\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi(x) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle\right]
\end{align*}
$$

- The first term gives

$$
\begin{align*}
& \left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| a_{\mathrm{out}}^{\dagger}\left(\vec{q}_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle \sqrt{2 E_{q_{1}}}= \\
& \left.=\sum_{k=1}^{n} 2 E_{p_{k}}(2 \pi)^{3} \delta^{(3)}\left(\vec{p}_{k}-\vec{q}_{1}\right)\left\langle p_{1}, \ldots, \hat{p}_{k}, \ldots, p_{n}, \text { out }\right| q_{2}, \ldots, q_{r}, \text { in }\right\rangle \tag{2.56}
\end{align*}
$$

where $\hat{p}_{k}$ has to be taken out. This describes a process where one of the in- and outgoing states are identical and do not participate in scattering. Such an amplitude corresponds to a disconnected diagram and its computation reduces to computing an S-matrix element involving only $(r-1)$ in- and $(n-1)$ out-states (since one in- and out-state factor out). The second term in (2.55) gives

$$
\begin{align*}
& i \int \mathrm{~d}^{4} x Z^{-1 / 2} \partial_{0}\left[e^{-i q_{1} \cdot x} \partial_{0}\langle\ldots\rangle-\left(\partial_{0} e^{-i q_{1} \cdot x}\right)\langle\ldots\rangle\right] \\
& =i \int \mathrm{~d}^{4} x Z^{-1 / 2}\left[e^{-i q_{1} \cdot x} \partial_{0}^{2}\langle\ldots\rangle-\left(\partial_{0}^{2} e^{-i q_{1} \cdot x}\right)\langle\ldots\rangle\right] \tag{2.57}
\end{align*}
$$

because the cross-terms cancel each other. Now consider that

$$
\begin{equation*}
-\partial_{0}^{2} e^{-i q_{1} \cdot x}=\left(q_{1}^{0}\right)^{2} e^{-i q_{1} \cdot x}=\left(q_{1}^{2}+\vec{q}_{1}^{2}\right) e^{-i q_{1} \cdot x}=\left(m^{2}-\nabla^{2}\right) e^{-i q_{1} \cdot x} . \tag{2.58}
\end{equation*}
$$

With respect to the spatial variable we can integrate two times by parts,

$$
\begin{equation*}
\int \mathrm{d}^{4} x\left(m^{2}-\nabla^{2}\right) e^{-i q_{1} \cdot x}\langle\ldots\rangle=\int \mathrm{d}^{4} x e^{-i q_{1} \cdot x}\left(m^{2}-\nabla^{2}\right)\langle\ldots\rangle, \tag{2.59}
\end{equation*}
$$

since boundary terms at spatial infinity vanish. To justify this recall that the momentum eigenstates are to be thought of as convoluted with a wavefunction profile as in $\left|f_{\text {in }}\right\rangle=\int \mathrm{d}^{3} p f(\vec{p})\left|p_{\text {in }}\right\rangle$ so that the full states are really localised in space. Altogether this gives

$$
\begin{align*}
& \left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| q_{1}, \ldots, q_{r}, \text { in }\right\rangle= \\
& \left.\sum_{k=1}^{n} 2 E_{p_{k}}(2 \pi)^{3} \delta^{(3)}\left(\vec{p}_{k}-\vec{q}_{1}\right)\left\langle p_{1}, \ldots, \hat{p}_{k}, \ldots, p_{n}, \text { out }\right| q_{2}, \ldots, q_{r}, \text { in }\right\rangle  \tag{2.60}\\
& \left.+i Z^{-1 / 2} \int \mathrm{~d}^{4} x_{1} e^{-i q_{1} \cdot x_{1}}\left(\square_{1}+m^{2}\right)\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi\left(x_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle .
\end{align*}
$$

- Now we repeat this for all remaining states. First consider replacing $\left\langle p_{1}\right|$ by $\phi_{\text {out }}$ as follows:

$$
\begin{align*}
& \left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| \phi\left(x_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle \\
& \left.=\sqrt{2 E_{p_{1}}}\left\langle p_{2}, \ldots, p_{n}, \text { out }\right| a_{\text {out }}\left(p_{1}\right) \phi\left(x_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle  \tag{2.61}\\
& \left.=\lim _{y_{1}^{0} \rightarrow \infty} i Z^{-1 / 2} \int \mathrm{~d}^{3} y_{1} e^{i p_{1} \cdot y_{1}} \stackrel{\leftrightarrow}{\partial_{y_{1}^{0}}}\left\langle p_{2}, \ldots, p_{n}, \text { out }\right| \phi\left(y_{1}\right) \phi\left(x_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle .
\end{align*}
$$

We would like to repeat the previous logic and transform this into a sum of two terms, one of which being the disconnected term $\left\langle p_{2}, \ldots, p_{n}\right.$, out $\left.\phi\left(x_{1}\right) a_{\text {in }}\left(\vec{p}_{1}\right)\right| q_{2}, \ldots, q_{r}$, in $\rangle$. However, we need to be careful with the ordering of operators as we cannot simply commute $a_{\text {in }}\left(\vec{p}_{1}\right)$ through $\phi\left(x_{1}\right)$. This is where the the time-ordering symbol $T$ comes in. Namely, observe that for finite values of $x_{1}^{0}$

$$
\begin{align*}
& \left.\lim _{\substack{t_{f} \rightarrow \infty \\
t_{i} \rightarrow-\infty}} \int_{t_{i}}^{t_{f}} \mathrm{~d} y_{1}^{0}\left[\left.\frac{\partial}{\partial y_{1}^{0}} i Z^{-1 / 2} \int \mathrm{~d}^{3} y_{1} e^{i p_{1} \cdot y_{1}} \stackrel{\leftrightarrow}{\partial_{y_{1}^{0}}^{0}}\left\langle p_{2}, \ldots, p_{n}, \text { out }\right| T \phi\left(y_{1}\right) \phi\left(x_{1}\right) \right\rvert\, q_{2}, \ldots, q_{r}, \text { in }\right\rangle\right] \\
& \left.=\lim _{y_{1}^{0} \rightarrow \infty}\left[i Z^{-1 / 2} \int \mathrm{~d}^{3} y_{1} e^{i p_{1} \cdot y_{1}} \stackrel{\leftrightarrow}{\partial_{y_{1}^{0}}}\left\langle p_{2}, \ldots, p_{n}, \text { out }\right| \phi\left(y_{1}\right) \phi\left(x_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle\right]  \tag{2.62}\\
& -\underbrace{-\sqrt{2 E_{p_{1}}} \rightarrow \text { disconnected term }}_{\left.\equiv\left\langle p_{2}, \ldots, \text { out }\right| \phi\left(x_{1}\right) a_{\text {in }}\left(\overrightarrow{p_{1}}\right) \mid q_{2}, \ldots, \text { in }\right\rangle} \lim _{y_{1} \rightarrow-\infty}\left[i Z^{-1 / 2} \int \mathrm{~d}^{3} y_{1} e^{i p_{1} \cdot y_{1}} \stackrel{\leftrightarrow}{\partial_{y_{1}^{0}}}\left\langle p_{2}, \ldots, p_{n}, \text { out }\right| \phi\left(x_{1}\right) \phi\left(y_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle] . .
\end{align*}
$$

Note the different ordering of $\phi\left(y_{1}\right) \phi\left(x_{1}\right)$ in both terms. Indeed since the limit $x_{1}^{0} \rightarrow \pm \infty$ appears outside of the correlator in eq. (2.60), the time-ordering symbol precisely yields these orderings as $y_{1}^{0} \rightarrow \pm \infty$. The term on the lefthand-side of (2.62) is, as before,

$$
\begin{equation*}
\left.i Z^{-1 / 2} \int \mathrm{~d}^{4} y_{1} e^{i p_{1} \cdot y_{1}}\left(\square_{y_{1}}+m^{2}\right)\left\langle p_{2}, \ldots, p_{n}, \text { out }\right| T \phi\left(y_{1}\right) \phi\left(x_{1}\right) \mid q_{2}, \ldots, q_{r}, \text { in }\right\rangle \tag{2.63}
\end{equation*}
$$

- This can be repeated for all in- and out states to get the Lehmann-Symanzik-Zimmermann reduction formula

$$
\begin{array}{|l}
\left.\left.\left\langle p_{1}, \ldots, p_{n}, \text { out }\right| q_{1}, \ldots, q_{r}, \text { in }\right\rangle \equiv\left\langle p_{1}, \ldots, p_{n}, \text { in }\right| S \mid q_{1}, \ldots, q_{r}, \text { in }\right\rangle \\
=\left(\sum \text { disconnected terms }\right)+ \\
+\left(i Z^{-1 / 2}\right)^{n+r} \int \mathrm{~d}^{4} y_{1} \ldots \mathrm{~d}^{4} y_{n} \int \mathrm{~d}^{4} x_{1} \ldots \mathrm{~d}^{4} x_{r}  \tag{2.64}\\
\times e^{i\left(\sum_{k=1}^{n} p_{k} \cdot y_{k}-\sum_{l=1}^{r} q_{l} \cdot x_{l}\right)} \\
\times\left(\square_{y_{1}}+m^{2}\right) \ldots\left(\square_{x_{1}}+m^{2}\right) \ldots\langle\Omega| T \phi\left(y_{1}\right) \ldots \phi\left(y_{n}\right) \phi\left(x_{1}\right) \ldots \phi\left(x_{r}\right)|\Omega\rangle .
\end{array}
$$

This formula reduces the computation of the $S$-matrix to the computation of time-ordered correlation functions of the full interacting theory.

- In terms of the Fourier transformed quantities it reads as follows. First note that

$$
\begin{equation*}
\left(\square_{y}+m^{2}\right) \phi(y)=\int \frac{\mathrm{d}^{4} \tilde{p}}{(2 \pi)^{4}}\left(-\tilde{p}^{2}+m^{2}\right) e^{-i \tilde{p} \cdot y} \tilde{\phi}(\tilde{p}) . \tag{2.65}
\end{equation*}
$$

We can plug this into (2.64), perform the integrals $\int d^{4} y_{k} e^{i\left(p_{k}-\tilde{p}_{k}\right) y_{k}}=(2 \pi)^{4} \delta^{(4)}\left(p_{k}-\tilde{p}_{k}\right)$ (and similarly for $x_{l}$, where we define accordingly $\left.\left(\square_{y}+m^{2}\right) \phi(x)=\int \frac{\mathrm{d}^{2} \tilde{q}}{(2 \pi)^{4}}\left(-\tilde{q}^{2}+m^{2}\right) e^{i \tilde{q} \cdot x} \tilde{\phi}(\tilde{q})\right)$ and arrive at

$$
\begin{align*}
\left.\left\langle p_{1}, \ldots p_{n}\right| S\left|q_{1}, \ldots, q_{r}\right\rangle\right|_{\text {connected }}= & \left(i Z^{-1 / 2}\right)^{n+r} \prod_{k=1}^{n}\left(-p_{k}^{2}+m^{2}\right) \prod_{l=1}^{r}\left(-q_{l}^{2}+m^{2}\right)  \tag{2.66}\\
& \times\langle\Omega| T \tilde{\phi}\left(p_{1}\right) \ldots \tilde{\phi}\left(p_{n}\right) \tilde{\phi}\left(q_{1}\right) \ldots \tilde{\phi}\left(q_{r}\right)|\Omega\rangle .
\end{align*}
$$

Note that the $p_{1}, \ldots, p_{n}$ and $q_{1}, \ldots q_{r}$ which appear on both sides of this equation are on-shell since these correspond to the physical 4-momenta of the out- and incoming 1-particles states. Therefore $p_{k}^{2}-m^{2}=0=q_{l}^{2}-m^{2}$. In order for $\left.\left\langle p_{1}, \ldots p_{n}\right| S\left|q_{1}, \ldots, q_{r}\right\rangle\right|_{\text {connected }}$ to be non-zero, the correlation function appearing on the right must therefore have a suitable pole structure such as to cancel precisely the kinematic factors $\prod_{k=1}^{n}\left(-p_{k}^{2}+m^{2}\right) \prod_{l=1}^{r}\left(-q_{l}^{2}+m^{2}\right)$.
In fact, as will be confirmed by explicit computation, $\langle\Omega| T \tilde{\phi}\left(p_{1}\right) \ldots \tilde{\phi}\left(p_{n}\right) \tilde{\phi}\left(q_{1}\right) \ldots \tilde{\phi}\left(q_{r}\right)|\Omega\rangle$ will in general be a sum of terms with different poles in the momenta. Only the term with the pole structure given precisely by $\prod_{k=1}^{n} \frac{1}{p_{k}^{2}-m^{2}} \prod_{l=1}^{n} \frac{1}{q_{l}^{2}-m^{2}}$ contributes to the connected S-matrix element $\left.\left\langle p_{1}, \ldots p_{n}\right| S\left|q_{1}, \ldots, q_{r}\right\rangle\right|_{\text {connected }}$. The terms with fewer poles will contribute at best to disconnected scattering processes. On the other hand, since the $S$-matrix, being a QM probability amplitude, is non-singular, the correlation functions cannot have more poles than what is cancelled by the kinematic factors on the right, and this prediction of the LSZ-formula will indeed be confirmed in explicit computations.

Thus

$$
\begin{align*}
& \prod_{k=1}^{n} \int \mathrm{~d}^{4} y_{k} e^{i p_{k} \cdot y_{k}} \prod_{l=1}^{r} \int \mathrm{~d}^{4} x_{l} e^{-i q_{l} \cdot x_{l}}\langle\Omega| T \prod_{k} \phi\left(y_{k}\right) \prod_{l} \phi\left(x_{l}\right)|\Omega\rangle  \tag{2.67}\\
& \quad=\left.\left(\prod_{k=1}^{n} \frac{i \sqrt{Z}}{p_{k}^{2}-m^{2}}\right)\left(\prod_{l=1}^{r} \frac{i \sqrt{Z}}{q_{l}^{2}-m^{2}}\right)\left\langle p_{1}, \ldots p_{n}\right| S\left|q_{1}, \ldots, q_{r}\right\rangle\right|_{\text {connected }}
\end{align*}
$$

We have arrived at a very precise prescription to compute the connected piece of the S-matrix element $\left.\left\langle p_{1}, \ldots p_{n}\right| S\left|q_{1}, \ldots, q_{r}\right\rangle\right|_{\text {connected }}:$

- Compute the Fourier transformation of the corresponding time-ordered $(n+r)$-correlation function and take all momenta $p_{k}$ and $q_{l}$ on-shell, i.e. $p_{k}^{0}=\sqrt{\vec{p}_{k}^{2}+m^{2}}$ and $q_{l}^{0}=\sqrt{\vec{p}_{l}^{2}+m^{2}}$.
- The result will be a sum of terms each a function of the momenta, which are distinguished by the structure of their poles.
- The connected $S$-matrix element (times $\left.(i \sqrt{Z})^{n+r}\right)$ ) is the residue with respect to $\prod_{k=1}^{n} \frac{1}{p_{k}^{2}-m^{2}} \prod_{l=1}^{r} \frac{1}{q_{l}^{2}-m^{2}}$.


### 2.5 Correlators in the interaction picture

To further evaluate the LSZ formula we need to compute the $(\mathrm{n}+\mathrm{r})$-correlation function

$$
\begin{equation*}
\langle\Omega| T \phi\left(y_{1}\right) \ldots \phi\left(y_{n}\right) \phi\left(x_{1}\right) \ldots \phi\left(x_{r}\right)|\Omega\rangle \tag{2.68}
\end{equation*}
$$

of the full interacting theory with Hamiltonian

$$
\begin{equation*}
H=H_{0}+H_{\mathrm{int}} . \tag{2.69}
\end{equation*}
$$

There are two ways how to do this, either by performing a path-integral computation, or by computing in the Interaction Picture. The path-integral formalism is reserved for the course QFT II. In the sequel we consider the latter approach. Our strategy is to reduce the computation of the full correlator to a calculation in terms of

- free-field creation/annihilation operators and
- the free-field vacuum $|0\rangle$.

This is achieved in the Interaction Picture ( $\equiv$ Dirac picture). Let $\phi(t, \vec{x})$ denote the Heisenberg Picture field of the full interacting theory and fix some reference time $t_{0}$. Then we define the Interaction Picture operators

$$
\begin{align*}
& \Phi_{I}(t, \vec{x})=e^{i H_{0}\left(t-t_{0}\right)} \phi\left(t_{0}, \vec{x}\right) e^{-i H_{0}\left(t-t_{0}\right)}  \tag{2.70}\\
& \Pi_{I}(t, \vec{x})=e^{i H_{0}\left(t-t_{0}\right)} \Pi\left(t_{0}, \vec{x}\right) e^{-i H_{0}\left(t-t_{0}\right)} .
\end{align*}
$$

The motivation behind this definition is that $\Phi_{I}(t, \vec{x})$ satisfies the free Klein-Gordon-equation

$$
\begin{equation*}
\left(\partial^{2}+m_{0}^{2}\right) \Phi_{I}(t, \vec{x})=0 \tag{2.71}
\end{equation*}
$$

with mass $m_{0}$ as in $H_{0}$. One can see this as follows: The defintion (2.70) implies that

$$
\begin{align*}
\partial_{t} \Phi_{I}(t, \vec{x}) & =i\left[H_{0}, \Phi_{I}(t, \vec{x})\right] \\
& =e^{i H_{0}\left(t-t_{0}\right)} \underbrace{i\left[H_{0}, \phi\left(t_{0}, \vec{x}\right)\right]}_{\Pi\left(t_{0}, \vec{x}\right)} e^{-i H_{0}\left(t-t_{0}\right)} . \tag{2.72}
\end{align*}
$$

Here we are using that $\Pi\left(t_{0}, \vec{x}\right)=i\left[H, \phi\left(t_{0}, \vec{x}\right)\right]=i\left[H_{0}, \phi\left(t_{0}, \vec{x}\right)\right]$ because $H$ and $H_{0}$ differ only by powers in $\phi(t, \vec{x}) .{ }^{4}$
Therefore

$$
\begin{equation*}
\partial_{t} \Phi_{I}(t, \vec{x})=\Pi_{I}(t, \vec{x}) \tag{2.73}
\end{equation*}
$$

and likewise

$$
\begin{align*}
\partial_{t}^{2} \Phi_{I}(t, \vec{x}) & =e^{i H_{0}\left(t-t_{0}\right)} \underbrace{i\left[H_{0}, \Pi\left(t_{0}, \vec{x}\right)\right]}_{\left(\nabla^{2}-m_{0}^{2}\right) \phi\left(t_{0}, \vec{x}\right)} e^{-i H_{0}\left(t-t_{0}\right)}  \tag{2.74}\\
& =\quad\left(\nabla^{2}-m_{0}^{2}\right) \Phi_{I}(t, \vec{x}) . \tag{2.75}
\end{align*}
$$

[^8]Here it is crucial to appreciate that $m_{0}$ appears because the commutator involves only $H_{0}$ and the computation of $i\left[H_{0}, \Pi\left(t_{0}, \vec{x}\right)\right]$ proceeds as in the free theory.
Equ. (2.71) implies that $\Phi_{I}(t, \vec{x})$ enjoys a free mode expansion of the form

$$
\begin{equation*}
\Phi_{I}(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(a_{I}(\vec{p}) e^{-i p \cdot x}+a_{I}^{\dagger}(\vec{p}) e^{i p \cdot x}\right), \quad p^{0}=E_{\vec{p}}=\vec{p}^{2}+m_{0}^{2} . \tag{2.76}
\end{equation*}
$$

Furthermore it is easy to see that the interaction picture fields and, as a result of (2.73), also the modes satisfy the free-field commutation relations

$$
\begin{equation*}
\left[\Phi_{I}(t, \vec{x}), \Pi_{I}(t, \vec{y})\right]=i \delta^{(3)}(\vec{x}-\vec{y}), \quad\left[a_{I}(\vec{p}), a_{I}^{\dagger}(\vec{q})\right]=(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{q}) . \tag{2.77}
\end{equation*}
$$

One then verifies that

$$
\begin{equation*}
\left[H_{0}, a_{I}(\vec{p})\right]=-E_{p} a_{I}(\vec{p}), \quad\left[H_{0}, a_{I}^{\dagger}(\vec{p})\right]=+E_{p} a_{I}^{\dagger}(\vec{p}) \tag{2.78}
\end{equation*}
$$

as in the free theory. This can be seen e.g. by noting that $H_{0}=e^{i H_{0}\left(t-t_{0}\right)} H_{0} e^{-i H_{0}\left(t-t_{0}\right)} \equiv\left(H_{0}\right)_{I}$ and therefore we can replace in $H_{0}$ all fields $\phi(x)$ by $\phi_{I}(x)$ so that all free field results carry over. Consequently by the same arguments as in the free theory there must exist a vacuum state annihilated by all $a_{I}(\vec{p})$ and by $H_{0}$. This identifies this state as the unique vacuum of the free theory, $H_{0}|0\rangle=0$, and therefore

$$
\begin{equation*}
a_{I}(\vec{p})|0\rangle=0 . \tag{2.79}
\end{equation*}
$$

At $t \neq t_{0}$, the Heisenberg Picture $\phi(t, \vec{x})$ and the Interaction Picture $\Phi_{I}(t, \vec{x})$ relate as

$$
\begin{align*}
\phi(t, \vec{x}) & =e^{i H\left(t-t_{0}\right)} \phi\left(t_{0}, \vec{x}\right) e^{-i H\left(t-t_{0}\right)} \\
& =e^{i H\left(t-t_{0}\right)} e^{-i H_{0}\left(t-t_{0}\right)} \Phi_{I}(t, \vec{x}) e^{i H_{0}\left(t-t_{0}\right)} e^{-i H\left(t-t_{0}\right)} \tag{2.80}
\end{align*}
$$

which yields

$$
\begin{equation*}
\phi(t, \vec{x})=U^{\dagger}\left(t, t_{0}\right) \Phi_{I}(t, \vec{x}) U\left(t, t_{0}\right), \tag{2.81}
\end{equation*}
$$

where

$$
\begin{equation*}
U\left(t, t_{0}\right)=e^{i H_{0}\left(t-t_{0}\right)} e^{-i H\left(t-t_{0}\right)} \tag{2.82}
\end{equation*}
$$

is the time-evolution operator. Note that, since $H$ and $H_{0}$ do not commute, this is not just $e^{-i H_{\text {int }}\left(t-t_{0}\right)}$. The logic is now to replace the Heisenberg Picture operators $\phi(x)$ in the correlator (2.68) by the Interaction Picture operators $\Phi_{I}(x)$ because these obey a free-mode expansion. To further evaluate the resulting expression, we first derive a useful expression for $U\left(t, t_{0}\right)$ and second establish a relation between the vacuum $|\Omega\rangle$ of the interacting theory as appearing in (2.68) and the free vacuum $|0\rangle$ on which the interaction picture modes act.

### 2.5.1 Time evolution

To compute $U\left(t, t_{0}\right)$ we note that the time-evolution operator satisfies

$$
\begin{equation*}
i \frac{\partial}{\partial t} U\left(t, t_{0}\right)=H_{I}(t) U\left(t, t_{0}\right) \tag{2.83}
\end{equation*}
$$

with $H_{I}(t)=e^{i H_{0}\left(t-t_{0}\right)} H_{\mathrm{int}} e^{-i H_{0}\left(t-t_{0}\right)}$ because

$$
\begin{align*}
i \partial_{t} e^{i H_{0}\left(t-t_{0}\right)} e^{-i H\left(t-t_{0}\right)} & =i\left[i H_{0} e^{i H_{0}\left(t-t_{0}\right)} e^{-i H\left(t-t_{0}\right)}+e^{i H_{0}\left(t-t_{0}\right)}(-i H) e^{-i H\left(t-t_{0}\right)}\right] \\
& =e^{i H_{0}\left(t-t_{0}\right)}\left(H-H_{0}\right) e^{-i H\left(t-t_{0}\right)} \\
& =e^{i H_{0}\left(t-t_{0}\right)} H_{\mathrm{int}} e^{-i H_{0}\left(t-t_{0}\right)} e^{i H_{0}\left(t-t_{0}\right)} e^{-i H\left(t-t_{0}\right)}  \tag{2.84}\\
& =H_{I}(t) U\left(t, t_{0}\right)
\end{align*}
$$

As in Quantum Mechanics we solve the differential equation (2.83)by rewriting it as an integral equation,

$$
\begin{equation*}
U\left(t, t_{0}\right)=\frac{1}{i} \int_{t_{0}}^{t} H_{I}\left(t^{\prime}\right) U\left(t^{\prime}, t_{0}\right) \mathrm{d} t^{\prime}+U\left(t_{0}, t_{0}\right) \tag{2.85}
\end{equation*}
$$

The latter can be solved iteratively with initial value $U^{(0)}\left(t, t_{0}\right)=\mathbb{1}$,

$$
\begin{align*}
& U^{(1)}\left(t, t_{0}\right)=\mathbb{1}+\frac{1}{i} \int_{t_{0}}^{t} \mathrm{~d} t_{1} H_{I}\left(t_{1}\right) \underbrace{U^{(0)}\left(t_{1}, t_{0}\right)}_{=\mathbb{1}}, \\
& U^{(2)}\left(t, t_{0}\right)=\mathbb{1}+\frac{1}{i} \int_{t_{0}}^{t} \mathrm{~d} t_{2} H_{I}\left(t_{2}\right) U^{(1)}\left(t_{2}, t_{0}\right), \tag{2.86}
\end{align*}
$$

and so on at each iteration. The exact solution is given by the $n$-th iteration for $n \rightarrow \infty$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} U^{(n)}\left(t, t_{0}\right) \equiv U\left(t, t_{0}\right) \tag{2.87}
\end{equation*}
$$

Indeed it is easy to check that this expression, which is just

$$
\begin{equation*}
U\left(t, t_{0}\right)=\mathbb{1}+\sum_{n=1}^{\infty}\left(\frac{1}{i}\right)^{n} \int_{t_{0}}^{t} \mathrm{~d} t_{1} \int_{t_{0}}^{t_{1}} \mathrm{~d} t_{2} \ldots \int_{t_{0}}^{t_{n-1}} \mathrm{~d} t_{n} H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right) \ldots H_{I}\left(t_{n}\right), \tag{2.88}
\end{equation*}
$$

solves equ. (2.83). Note that the $H_{I}\left(t_{i}\right)$ under the integral are time-ordered.
The solution can be simplified further. Consider e.g. the second term and observe that it can be rewritten as

$$
\begin{equation*}
\int_{t_{0}}^{t} \mathrm{~d} t_{1} \int_{t_{0}}^{t_{1}} \mathrm{~d} t_{2} H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right)=\frac{1}{2} \int_{t_{0}}^{t} \mathrm{~d} t_{1} \int_{t_{0}}^{t} \mathrm{~d} t_{2} T\left(H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right)\right) \tag{2.89}
\end{equation*}
$$

Indeed the right-hand side is

$$
\begin{equation*}
\underbrace{\frac{1}{2} \int_{t_{0}}^{t} \mathrm{~d} t_{1} \int_{t_{0}}^{t_{1}} \mathrm{~d} t_{2} H_{I}\left(t_{1}\right) H_{I}\left(t_{2}\right)}_{\text {a) } t_{1}>t_{2}}+\underbrace{\frac{1}{2} \int_{t_{0}}^{t} \mathrm{~d} t_{1} \int_{t_{1}}^{t} \mathrm{~d} t_{2} H_{I}\left(t_{2}\right) H_{I}\left(t_{1}\right)}_{\text {b) } t_{2}>t_{1}} \tag{2.90}
\end{equation*}
$$

and intgral $\mathbf{b}$ ) is in fact the same as integral a). To see this rotate the square in the $t_{1}-t_{2}$ plane over which we integrating by $90^{\circ}$, which gives

$$
\begin{equation*}
\text { b) } \left.=\frac{1}{2} \int_{t_{0}}^{t} \mathrm{~d} t_{2} \int_{t_{0}}^{t_{2}} \mathrm{~d} t_{1} H_{I}\left(t_{2}\right) H_{I}\left(t_{1}\right) \equiv \mathrm{a}\right) . \tag{2.91}
\end{equation*}
$$

With this reasoning the time-evolution operator is

$$
\begin{equation*}
U\left(t, t_{0}\right)=\sum_{n=0}^{\infty} \frac{(-i)^{n}}{n!} \int_{t_{0}}^{t} \mathrm{~d} t_{1} \int_{t_{0}}^{t} \mathrm{~d} t_{2} \ldots \int_{t_{0}}^{t} \mathrm{~d} t_{n} T H_{I}\left(t_{1}\right) H_{I}\left(t_{n}\right) . \tag{2.92}
\end{equation*}
$$

For the latter infinite series one introduces the notation

$$
\begin{equation*}
U\left(t, t_{0}\right)=T e^{-i \int_{t_{0}}^{t} \mathrm{~d}^{\prime} H_{l}\left(t^{\prime}\right)} \tag{2.93}
\end{equation*}
$$

From the series expression one can verify that the time-evolution operator has the following properties:

- $U^{\dagger}\left(t_{1}, t_{2}\right)=U^{-1}\left(t_{1}, t_{2}\right)=U\left(t_{2}, t_{1}\right)$,
- $U\left(t_{1}, t_{2}\right) U\left(t_{2}, t_{3}\right)=U\left(t_{1}, t_{3}\right)$ for $t_{1} \geq t_{2} \geq t_{3}$.

> | Note that here the |
| :--- |
| definition of $U$ is extended |
| to arbitrary values of the |
| second argument by |
| (2.93). |

### 2.5.2 From the interacting to the free vacuum

Having understood how to relate the full Heisenberg fields to the free interaction picture fields, we now try to set up a relation between the free vacuum $|0\rangle$ and the interacting vacuum $|\Omega\rangle$. Let $|n\rangle$ be an eigenstate of the full Hamiltonian, i.e.

$$
\begin{equation*}
H|n\rangle=E_{n}|n\rangle \text { with } H=H_{0}+H_{\mathrm{int}} \tag{2.94}
\end{equation*}
$$

Then the time-evolution of the free vacuum $|0\rangle$ is ${ }^{5}$

$$
\begin{align*}
e^{-i H T}|0\rangle & =e^{-i H T} \sum_{n}|n\rangle\langle n \mid 0\rangle \\
& =\sum_{n} e^{-i E_{n} T}|n\rangle\langle n \mid 0\rangle  \tag{2.95}\\
& =e^{-i E_{\Omega} T}|\Omega\rangle\langle\Omega \mid 0\rangle+\sum_{|n\rangle \neq|\Omega\rangle} e^{-i E_{n} T}|n\rangle\langle n \mid 0\rangle .
\end{align*}
$$

[^9]If $H_{0}|0\rangle=0$, then $H|\Omega\rangle=E_{\Omega}|\Omega\rangle$ with $E_{\Omega} \neq 0$, because we now compare the vacuum energy of 2 theories or, put differently, the counter-term $V_{0}$ in the Lagrangian has already been used to set $E_{0}=0$ so that we are stuck with $E_{\Omega}$, whatever it is. Be that as it may, we have $E_{n}>E_{\Omega} \forall|n\rangle \neq|\Omega\rangle$. So if we formally take the limit $T \rightarrow \infty(1-i \epsilon)$, then $e^{-i E_{n} T}$ is stronger suppressed and only the vacuum $|\Omega\rangle$ survives in ${ }^{6}$

$$
\begin{equation*}
\lim _{T \rightarrow \infty(1-i \epsilon)} e^{-i H T}|0\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} e^{-i E_{\Omega} T}\langle\Omega \mid 0\rangle|\Omega\rangle . \tag{2.96}
\end{equation*}
$$

Solving this for $|\Omega\rangle$ yields

$$
\begin{equation*}
|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{e^{-i H T}|0\rangle}{e^{-i E_{\Omega} T}\langle\Omega \mid 0\rangle} \tag{2.97}
\end{equation*}
$$

Note that we are assuming here that $\langle\Omega||0\rangle \neq 0$, which is guaranteed at least for small perturbations $H_{\text {int }}$ in $H$.
To bring this to a form involving the time-evolution operator we shift $\left.T \rightarrow T+t_{0}=t_{0}-(-T)\right)$ and write for $e^{-i H\left(T+t_{0}\right)}|0\rangle$

$$
\begin{equation*}
e^{-i H\left(t_{0}-(-T)\right)} \underbrace{e^{-i H_{0}\left(-T-t_{0}\right)}|0\rangle}_{=|0\rangle \text { since } H_{0}|0\rangle=0}=U\left(t_{0},-T\right)|0\rangle, \tag{2.98}
\end{equation*}
$$

because

$$
\begin{align*}
U\left(t_{0},-T\right)=U\left(-T, t_{0}\right)^{-1} & =\left[e^{i H_{0}\left(-T-t_{0}\right)} e^{-i H\left(-T-t_{0}\right)}\right]^{-1}  \tag{2.99}\\
& =e^{-i H\left(T+t_{0}\right)} e^{-i H_{0}\left(-T-t_{0}\right)} .
\end{align*}
$$

So the vacuum is

$$
\begin{equation*}
|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)}\left(e^{-i E_{\Omega}\left(t_{0}-(-T)\right)}\langle\Omega \mid 0\rangle\right)^{-1} U\left(t_{0},-T\right)|0\rangle . \tag{2.100}
\end{equation*}
$$

Likewise, starting from $\lim _{T \rightarrow \infty(1-i \epsilon)}\langle 0| e^{-i H T}$ one gets

$$
\begin{equation*}
\langle\Omega|=\lim _{T \rightarrow \infty(1-i \epsilon)}\langle 0| U\left(T, t_{0}\right)\left(e^{-i E_{\Omega}\left(T-t_{0}\right)}\langle 0 \mid \Omega\rangle\right)^{-1} . \tag{2.101}
\end{equation*}
$$

Finally we can compute $\langle\Omega| T \phi(x) \phi(y)|\Omega\rangle$. Suppose first that $x^{0} \geq y^{0} \geq t_{0}$ : Then $\langle\Omega| T \phi(x) \phi(y)|\Omega\rangle$ becomes

$$
\begin{align*}
\langle\Omega| \phi(x) \phi(y)|\Omega\rangle= & \lim _{T \rightarrow \infty(1-i \epsilon)}\left(e^{-i E_{\Omega}\left(T-t_{0}\right)}\langle 0 \mid \Omega\rangle\right)^{-1} \\
& \times\langle 0| U\left(T, t_{0}\right) \underbrace{U\left(x^{0}, t_{0}\right)^{\dagger} \Phi_{I}(x) U\left(x^{0}, t_{0}\right)}_{\phi(x)} \\
& \times \underbrace{U\left(y^{0}, t_{0}\right)^{\dagger} \Phi_{I}(x) U\left(y^{0}, t_{0}\right)}_{\phi(y)} U\left(t_{0},-T\right)|0\rangle\left(e^{-i E_{\Omega}\left(t_{0}-(-T)\right)}\langle\Omega \mid 0\rangle\right)^{-1}  \tag{2.102}\\
= & \lim _{T \rightarrow \infty(1-i \epsilon)}\left(|\langle 0 \mid \Omega\rangle|^{2} e^{-i E_{\Omega} 2 T}\right)^{-1} \\
& \times\langle 0| U\left(T, x_{0}\right) \Phi_{I}(x) U\left(x^{0}, y^{0}\right) \Phi_{I}(y) U\left(y^{0},-T\right)|0\rangle .
\end{align*}
$$

[^10]Note that the contraction $U\left(x^{0}, t_{0}\right) U\left(t_{0}, y^{0}\right)=U\left(x^{0}, y^{0}\right)$ only works because of time-ordering. We can eliminate the constant prefactor by noting that

$$
\begin{equation*}
\mathbb{1}=\langle\Omega \mid \Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)}\left(|\langle 0 \mid \Omega\rangle|^{2} e^{-i E_{\Omega}(2 T)}\right)^{-1}\langle 0| U\left(T, t_{0}\right) U\left(t_{0},-T\right)|0\rangle . \tag{2.103}
\end{equation*}
$$

For $x^{0} \geq y^{0} \geq t_{0}$ we arrive at

$$
\begin{equation*}
\langle\Omega| \phi(x) \phi(y)|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{\langle 0| U\left(T, x^{0}\right) \Phi_{I}(x) U\left(x^{0}, y^{0}\right) \Phi_{I}(y) U\left(y^{0},-T\right)|0\rangle}{\langle 0| U(T,-T)|0\rangle} . \tag{2.104}
\end{equation*}
$$

The nominator is

$$
\begin{equation*}
\langle 0| T \Phi_{I}(x) \Phi_{I}(y) \underbrace{U\left(T, x^{0}\right) U\left(x^{0}, y^{0}\right) U\left(y^{0},-T\right)}_{U(T,-T)}|0\rangle, \tag{2.105}
\end{equation*}
$$

where the time-ordering symbol takes care of order. Similar conclusions are obtained for $x^{0} \leq y^{0}$. Therefore altogether

$$
\begin{equation*}
\langle\Omega| T \phi(x) \phi(y)|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{\langle 0| T\left(\Phi_{I}(x) \Phi_{I}(y) e^{-i} \int_{-T}^{T} \mathrm{~d} t H_{I}(t)\right.}{}|0\rangle . \tag{2.106}
\end{equation*}
$$

This same reasoning goes through for higher $n$-point correlators. Our master formula for computing correlation function becomes

$$
\begin{equation*}
\langle\Omega| T \prod_{i} \phi\left(x_{i}\right)|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{\langle 0| T \prod_{i} \Phi_{I}\left(x_{i}\right) e^{-i} \int_{-T}^{T} \mathrm{~d} t H_{l}(t)}{\langle 0\rangle} \frac{\langle 0| T e^{-i} \int_{-T}^{T} \mathrm{~d} H_{l}(t)}{}|0\rangle . \tag{2.107}
\end{equation*}
$$

This formula can be applied to concrete interactions, i.e.

$$
\begin{equation*}
\mathcal{L}=\mathcal{L}_{0}+\mathcal{L}_{\mathrm{int}}, \text { e.g. } \mathcal{L}_{\mathrm{int}}=-\frac{\lambda}{4!} \phi^{4}(x), \tag{2.108}
\end{equation*}
$$

where $\lambda$ has mass-dimension 0 . The interaction picture Hamiltonian is then

$$
\begin{equation*}
H_{I}=\int \mathrm{d}^{3} x \frac{\lambda}{4!} \Phi_{I}(t, \vec{x})^{4} \tag{2.109}
\end{equation*}
$$

and we can expand the time-evolution perturbatively order by order in $\lambda$ if $\lambda \ll 1$ by expanding the exponential

$$
\begin{equation*}
\lim _{T \rightarrow \infty(1-i \epsilon)} e^{-i \int_{-T}^{T} \mathrm{~d} t H_{l}(t)}=\lim _{T \rightarrow \infty(1-i \epsilon)} e^{-i \int \mathrm{~d}^{4} x_{4}^{\lambda} \Phi_{I}^{4}(x)} . \tag{2.110}
\end{equation*}
$$

As we will see the ( $1-i \epsilon$ ) prescription for the boundaries of the integral will pose no problems (cf. discussion after (2.133)).
There are basically two remaining problems:

- Perform a systematic evaluation of $\langle 0| T \prod_{i} \Phi_{I}\left(x_{i}\right)|0\rangle$ and
- deal with the denominator.

We will solve these problems by exploiting the action of the creation and annihilation operators on the vacuum.

### 2.6 Wick's theorem

To compute an expression of the form $\langle 0| T \prod_{i} \Phi_{I}\left(x_{i}\right)|0\rangle$ we decompose $\Phi_{I}$ into free modes,

$$
\begin{equation*}
\Phi_{I}=\underbrace{\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}} a_{I}(\vec{p}) e^{-i p \cdot x}}_{=: \Phi_{I}^{+}(x)}+\underbrace{\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}} a_{I}^{\dagger}(\vec{p}) e^{i p \cdot x}}_{=: \Phi_{I}^{-}(x)} . \tag{2.111}
\end{equation*}
$$

Now we want to commute all $\Phi_{I}^{-}(x)$ to the left of all $\Phi_{I}^{+}(x)$ to use

$$
\begin{equation*}
\Phi_{I}^{+}(x)|0\rangle=0=\langle 0| \Phi_{I}^{-}(x) \tag{2.112}
\end{equation*}
$$

To do so we definie normal-ordering:

An operator $O$ is normal-ordered if all creation/annihilation operators appear on the left/right. For such $O$ we write : $O$ :

Consider for example the operator:

$$
\begin{equation*}
: a^{\dagger}\left(\vec{p}_{1}\right) a\left(\vec{p}_{2}\right) a^{\dagger}\left(\vec{p}_{3}\right):=a^{\dagger}\left(\vec{p}_{1}\right) a^{\dagger}\left(\vec{p}_{3}\right) a\left(\vec{p}_{2}\right)=a^{\dagger}\left(\vec{p}_{3}\right) a^{\dagger}\left(\vec{p}_{1}\right) a\left(\vec{p}_{2}\right) \tag{2.113}
\end{equation*}
$$

It is obvious that

$$
\begin{equation*}
\langle 0|: O:|0\rangle=0 \tag{2.114}
\end{equation*}
$$

for every non-trivial operator $O \neq c \mathbb{1}$ for $c \in \mathbb{C}$. We begin with $\langle 0| T \Phi_{I}(x) \Phi_{I}(y)|0\rangle$ and drop the subscript " $I$ " from now on. There are two cases to consider, either $x^{0} \geq y^{0}$ or $y^{0} \geq x^{0}$. If $x^{0} \geq y^{0}$, then $T \Phi(x) \Phi(y)=\Phi(x) \Phi(y)$ and

$$
\begin{equation*}
\Phi(x) \Phi(y)=\Phi^{-}(x) \Phi^{-}(y)+\Phi^{-}(x) \Phi^{+}(y)+\Phi^{+}(x) \Phi^{+}(y)+\Phi^{+}(x) \Phi^{-}(y) \tag{2.115}
\end{equation*}
$$

The first three terms are already normal-ordered and the last term can be put in normal-ordered form by commuting the fields through each other,

$$
\begin{equation*}
\Phi(x) \Phi(y)=: \Phi(x) \Phi(y):+\left[\Phi^{+}(x), \Phi^{-}(y)\right] \tag{2.116}
\end{equation*}
$$

Similar expressions follow for $x^{0} \leq y^{0}$, so altogether

$$
\begin{align*}
T(\Phi(x) \Phi(y))=: \Phi(x) \Phi(y): & +\Theta\left(x^{0}-y^{0}\right)\left[\Phi^{+}(x), \Phi^{-}(y)\right]  \tag{2.117}\\
& +\Theta\left(y^{0}-x^{0}\right)\left[\Phi^{+}(y), \Phi^{-}(x)\right]
\end{align*}
$$

We notice that the last two terms are a $\mathbb{C}$-number $c$ and with

$$
\begin{equation*}
\langle 0| T \Phi(x) \Phi(y)|0\rangle=\underbrace{\langle 0|: \Phi(x) \Phi(y):|0\rangle}_{\equiv 0}+\langle 0| c|0\rangle=c \equiv D_{F}^{(0)}(x-y) \tag{2.118}
\end{equation*}
$$

we find that

$$
\begin{equation*}
T(\Phi(x) \Phi(y))=: \Phi(x) \Phi(y):+D_{F}^{(0)}(x-y) . \tag{2.119}
\end{equation*}
$$

We use the notation $D_{F}^{(0)}(x-y)$ for the free Feynman propagator in the sequel to emphasize that this object is defined in terms of the free mass parameter $m_{0}$ appearing in the Lagrangian. If we define the contraction

$$
\begin{equation*}
\widetilde{\Phi(x) \Phi(y)}=D_{F}^{(0)}(x-y)=D_{F}^{(0)}(y-x) \tag{2.120}
\end{equation*}
$$

we can write this as

$$
\begin{equation*}
T(\Phi(x) \Phi(y))=: \Phi(x) \Phi(y)+\widetilde{\Phi(x) \Phi(y)}: \tag{2.121}
\end{equation*}
$$

Note that by definition : $c:=c$ for a C-number, which explains the notation : $\overline{\Phi(x) \Phi(y)}::^{7}$ This generalises to higher products. E.g. one can show by directcomputation that

$$
\begin{aligned}
& T\left(\Phi\left(x_{1}\right) \Phi\left(x_{2}\right) \Phi\left(x_{3}\right)\right) \equiv T\left(\Phi_{1} \Phi_{2} \Phi_{3}\right) \\
& =: \Phi_{1} \Phi_{2} \Phi_{3}:+: \Phi_{1} \widetilde{\Phi_{2} \Phi_{3}}:+: \widetilde{\Phi_{1} \Phi_{2} \Phi_{3}:+: \widetilde{\Phi_{1} \Phi_{2} \Phi_{3}}: ~} \\
& =: \Phi_{1}: D_{F}^{(0)}\left(x_{2}-x_{3}\right)+: \Phi_{3}: D_{F}^{(0)}\left(x_{1}-x_{2}\right)+: \Phi_{2}: D_{F}^{(0)}\left(x_{1}-x_{3}\right) \text {. }
\end{aligned}
$$

Wick's theorem generalises this for $N$ fields,

$$
\begin{align*}
T\left(\Phi_{1} \ldots \Phi_{N}\right) & =: \Phi_{1} \ldots \Phi_{N}: \\
& +\sum_{1 \leq i<j \leq N}: \Phi_{1} \ldots \Phi_{i} \ldots \Phi_{j} \ldots \Phi_{N}: \\
& +\sum_{1 \leq i<k<j<l \leq N}\left(: \Phi _ { 1 } \ldots \Phi _ { i } \ldots \Phi _ { k } \ldots \longdiv { \Phi _ { j \ldots } \ldots \Phi _ { l } \ldots \Phi _ { N } : }\right.  \tag{2.123}\\
& \quad+2 \text { more contractions }) \\
& +\ldots
\end{align*}
$$

Indeed one can prove this by induction and summarize Wick's theorem as

$$
\begin{equation*}
T\left(\Phi_{1} \ldots \Phi_{N}\right)=: \Phi_{1} \ldots \Phi_{N}:+: \text { all contractions of distinct pairs : . } \tag{2.124}
\end{equation*}
$$

The proof is reserved to Examples Sheet 6.
Wick's theorem together with (2.114) has two important consequences:

- $\langle 0| T \Phi_{1} \ldots \Phi_{2 N+1}|0\rangle=0$, because there is always an odd number of normal-ordered fields remaining and
- $\langle 0| T \Phi_{1} \ldots \Phi_{2 N}|0\rangle=D_{F}^{(0)}\left(x_{1}-x_{2}\right) D_{F}^{(0)}\left(x_{3}-x_{4}\right) \ldots D_{F}^{(0)}\left(x_{2 N-1}-x_{2 N}\right)+\sum$ (all other contractions) .

[^11]
### 2.7 Feynman diagrams

There exists a practical graphical representation of the systematics of contractions in terms of Feynman diagrams. Consider for instance

$$
\begin{align*}
\langle 0| T\left(\Phi_{1} \Phi_{2} \Phi_{3} \Phi_{4}\right)|0\rangle & =D_{F}^{(0)}\left(x_{1}-x_{2}\right) D_{F}^{(0)}\left(x_{3}-x_{4}\right) \\
& +D_{F}^{(0)}\left(x_{1}-x_{3}\right) D_{F}^{(0)}\left(x_{2}-x_{4}\right)+D_{F}^{(0)}\left(x_{1}-x_{4}\right) D_{F}^{(0)}\left(x_{2}-x_{3}\right) . \tag{2.125}
\end{align*}
$$

The result can be translated into a Feynman diagram as follows:

- First, draw 1 point for all $x_{i}$ in $\Phi\left(x_{i}\right)$ and connect these by lines in all possible ways. This gives, in the present case, three distinct diagrams displayed in Figure 2.2. Note that these are distinct because we distinguish between $x_{1}, x_{2}, x_{3}, x_{4}$.


Figure 2.2: Possible connections between all $x_{i}$.

- Then, to each line between $i$ and $j$ we attach a factor of $D_{F}^{(0)}\left(x_{i}-x_{j}\right)$.
- We then multiply all $D_{F}^{(0)}\left(x_{i}-x_{j}\right)$ and sum over all distinct diagrams.

In interacting theories also contractions involving fields at the same spacetime point occur. For example consider the propagator in $\Phi^{4}$-theory up to first order in $\lambda$ :

$$
\begin{align*}
& \langle 0| T \Phi(x) \Phi(y) e^{-i \frac{\lambda}{4!} \int \mathrm{d}^{4} z \Phi^{4}(z)}|0\rangle= \\
& \quad\langle 0| T \Phi(x) \Phi(y)|0\rangle+\langle 0| T \Phi(x) \Phi(x) \int \mathrm{d}^{4} z\left(-\frac{i \lambda}{4!}\right) \Phi^{4}(z)|0\rangle+O\left(\lambda^{2}\right) . \tag{2.126}
\end{align*}
$$

The $O\left(\lambda^{0}\right)$-term corresponds to a straight line between $x$ and $y$, i.e. the propagation from $x$ to $y$, which is assigned a factor $D_{F}^{(0)}(x-y)$. The term to first order in $\lambda$ is

$$
\begin{align*}
& \frac{-i \lambda}{4!}\left[\int \mathrm{d}^{4} z \phi(x) \phi(y) \stackrel{\overparen{\Phi(z)} \sqrt{\Phi(z) \Phi(z)} \Phi(z)}{ }\right.  \tag{2.127}\\
& \left.+ \int \mathrm { d } ^ { 4 } z \Phi ( x ) \Phi ( y ) \Phi ( z ) \Phi ( z ) \longdiv { \Phi ( z ) \Phi ( z ) }\right] \tag{2.128}
\end{align*}
$$

This can be written as

$$
\begin{align*}
& \int \mathrm{d}^{4} z D_{F}^{(0)}(x-y)\binom{4}{2} \frac{1}{2} D_{F}^{(0)}(z-z) D_{F}^{(0)}(z-z) \\
& +\int \mathrm{d}^{4} z D_{F}^{(0)}(x-z) D_{F}^{(0)}(y-z) D_{F}^{(0)}(z-z) \cdot 4 \cdot 3 . \tag{2.129}
\end{align*}
$$

The graphical representation of these two different contributions is given in Figure 2.3. Note the appearance of the point $z$, which is integrated over. The first term is proportional $-i \lambda \frac{1}{8}$, while the second term has a factor of $-i \lambda \frac{1}{2}$.


Figure 2.3: First and second term of $O\left(\lambda^{1}\right)$.

To give an interpretation of these combinatorial factors of $1 / 8$ and $1 / 2$, we observe that they coincide with 1 over the order ( $\equiv$ number of elements) of the symmetry group of the diagrams. These are called the symmetry factors. The symmetry factors are given by the number of ways one can exchange components of the diagram without changing the diagram, where by components we mean either the two ends of a line starting and ending on the same point, entire lines between points or internal points (vertices - see below).
For the two diagrams these are in turn:

- The single loop at $z$ is symmetric under the exchange of the out- and in-going end of the loop, but the direction must be consistent (i.e. one out- and one ingoing arrow). Therefore the symmetry group is $G=\mathbb{Z}_{2}$ and its order is $|G|=2$.
- The double loop is separately symmetric under exchange of the out- and in-going end of both loops, which gives a $\mathbb{Z}_{2} \times \mathbb{Z}_{2}$ as well as under exchange of the two loops as a whole. Therefore the symmetry group is $G=\mathbb{Z}_{2} \times \mathbb{Z}_{2} \times \mathbb{Z}_{2}$ and its order is $|G|=8$.

One can show in general that if a Feynman diagram with symmetry group $G$ always carries a combinatorial factor $\frac{1}{|G|}$. This is because if some symmetry remains the factor of $\frac{1}{n!}$ in the interaction term $-\frac{\lambda}{n!} \phi^{n}$ is only partially cancelled by counting the various contractions that yield the same diagram. In the above examples the symmetry group was due to the ambiguity in assigning arrows to the lines at a vertex. In addition, if $k$ vertices are identical, the symmetry group includes the group of permutations of these $k$ vertices of order $\left|S_{k}\right|=k!$.

Let us introduce some jargon:

- The points $x, y$ associated to $\Phi(x) \Phi(y)$ in the correlator, which are not integrated over, are called external points.
- All other points - here z - over which is integrated over are internal points.
- For a $\Phi^{n}$ interaction, at an internal point $n$ lines meet. Such points are called vertices.


### 2.7.1 Position space Feynman-rules

The translation of Wick's theorem into Feynman-graphs as exemplified above can be proven to hold generally (the only non-trivial aspect being the symmetry factors). We summarize this in terms of the position space Feynman-rules for the computation of

$$
\begin{equation*}
\langle 0| T \prod_{i=1}^{m} \Phi\left(x_{i}\right) e^{-i \frac{\lambda}{n!} \int d^{4} z \phi^{n}(z)}|0\rangle \tag{2.130}
\end{equation*}
$$

at order $\lambda^{k}$, where we are assuming that all $x_{i}$ are distinct:

- Draw one external point for all $x_{i}$ and $k$ internal points $z_{j}$.
- Connect the points by lines such that
- to each external point $x_{i} 1$ line is attached,
- to each internal point $z_{j} n$ lines are attached.
- To each line between points $y_{i}$ and $y_{j}$ (both external and internal) we associate a free propagator

$$
\begin{equation*}
D_{F}^{(0)}\left(y_{i}-y_{j}\right)=D_{F}^{(0)}\left(y_{j}-y_{i}\right), \tag{2.131}
\end{equation*}
$$

where we stress again that the superscript reminds us to take the Feynman propagator of the free theory.

- To each vertex associate a factor $-i \lambda \int \mathrm{~d}^{4} z_{j}$.
- To each external point associate a 1.
- Multiply all factors, Feynman propagators etc. and divide by the symmetry factor of the diagram.
- Then sum up all distinct such Feynman diagrams.


### 2.7.2 Momentum space Feynman-rules

More practical are the momentum space Feynman-rules: To each free propagator

$$
\begin{equation*}
D_{F}^{(0)}(x-y)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i}{p^{2}-m_{0}^{2}+i \epsilon} e^{-i p \cdot(x-y)} \tag{2.132}
\end{equation*}
$$

associate an arrow and momentum $p$ in that direction. This means an arrow from $x$ to $y$ corresponds to the propagator $D_{F}^{(0)}(x-y)$. Obviously the direction is arbitrary because $D_{F}^{(0)}(x-y)=D_{F}^{(0)}(y-x)$,
but let us make one such choice for each line. We can then perform the integral $\int d^{4} z$ at each vertex explicitly. E.g. for the $\Phi^{4}$-theory 4 lines meet at a vertex, and if, say, the momenta $p_{1}$ and $p_{2}$ point into that vertex and $p_{3}, p_{4}$ point out of it, this integral yields

$$
\begin{equation*}
\int \mathrm{d}^{4} z e^{i p_{1} \cdot z} e^{i p_{2} \cdot z} e^{-i p_{3} \cdot z} e^{-i p_{4} \cdot z}=(2 \pi)^{4} \delta^{(4)}\left(p_{1}+p_{2}-p_{3}-p_{4}\right) . \tag{2.133}
\end{equation*}
$$

This implements momentum conservation at each vertex. The delta-function will effectively eliminate one integral $\int \mathrm{d}^{4} p_{i}$ per vertex. ${ }^{8}$
Note at this point that the formal prescription that we integrate $\lim _{T \rightarrow \infty(1-i \epsilon)} \int_{-T}^{T}$ in (2.107) translates into a corresponding prescription for $\int \mathrm{d}^{4} z$. We should therefore view this integral as a complex contour integral. From the fact that such integrals can be deformed in the complex plane as long as no poles are hit it follows that the $(1-i \epsilon)$ prescription does not affect the result.
The momentum-space Feynman rules ${ }^{9}$ for the computation of the $n$-point function

$$
\begin{equation*}
\langle 0| T \prod_{i=1}^{m} \Phi\left(x_{i}\right) e^{-i \frac{\lambda}{n!} \int d^{4} z \phi^{n}(z)}|0\rangle \tag{2.134}
\end{equation*}
$$

at order $\lambda^{k}$ are:

- Draw one external point for all $x_{i}$ and $k$ internal points $z_{j}$.
- Connect the points by lines such that
- to each external point $x_{i} 1$ line is attached,
- to each internal point $z_{j} n$ lines are attached.
- To each line between points $y_{i}$ and $y_{j}$ (both external and internal) we associate a free propagator

$$
\begin{equation*}
D_{F}^{(0)}\left(y_{i}-y_{j}\right) \tag{2.135}
\end{equation*}
$$

with one choice of direction and to each such $D_{F}^{(0)}\left(y_{i}-y_{j}\right)$ we associate directed momentum $p$ from $y_{i}$ to $y_{j}$ and a factor

$$
\begin{equation*}
\frac{i}{p^{2}-m_{0}^{2}+i \epsilon} . \tag{2.136}
\end{equation*}
$$

- For each vertex we multiply a factor of $(2 \pi)^{4} \delta^{(4)}\left(\sum_{\text {ingoing }} p_{i}-\sum_{\text {outgoing }} p_{k}\right) \times(-i \lambda)$.
- For each external point we multiply a factor of $e^{-i p \cdot x}$ for momentum pointing out of the external point, or $e^{i p \cdot x}$ for momentum pointing into the point.

[^12]- Integrate over each momentum $\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}}$ and divide by the symmetry factor.
- Then sum up all distinct such Feynman diagrams.

For example, consider again the $\Phi^{4}$ theory to order $\lambda$, i.e.

$$
\begin{equation*}
\langle 0| T \prod_{k=1}^{4} \Phi_{k} e^{-i \frac{\lambda}{4!} \int \mathrm{d}^{4} x \Phi^{4}}|0\rangle \tag{2.137}
\end{equation*}
$$

$\lambda^{0}$ corresponds to the diagrams in Figure 2.2. For $\lambda^{1}$ we can have


Figure 2.4: Diagrams for $\lambda^{1}$.

Consider the first diagram at order $\lambda$, make a consistent choice of momenta and assign $e^{i k \cdot x}$ for in-going arrows at external points, $e^{-i k \cdot x}$ for out-going arrows and $-i \lambda \delta^{(4)}\left(k_{1}+k_{2}-k_{3}-k_{4}\right)$ at the vertex due to momentum consvervation. We end up with

$$
\begin{align*}
\left.\langle 0| T \prod_{k=1}^{4} \Phi\left(x_{k}\right)\left(-\frac{i \lambda}{4!} \int \mathrm{d}^{4} z \Phi^{4}(z)\right)|0\rangle\right|_{\text {Diagram } \times} \\
\quad=\prod_{j=1}^{4} \int \frac{\mathrm{~d}^{4} k_{j}}{(2 \pi)^{4}} e^{-i k_{1} \cdot x_{1}} e^{-i k_{2} \cdot x_{2}} e^{i k_{3} \cdot x_{3}} e^{i k_{4} \cdot x_{4}}  \tag{2.138}\\
\quad \times(-i \lambda)(2 \pi)^{4} \delta^{(4)}\left(k_{1}+k_{2}-k_{3}-k_{4}\right) \prod_{j=1}^{4} \frac{i}{k_{j}^{2}-m_{0}^{2}+i \epsilon} .
\end{align*}
$$

As we will see soon, in the context of the LSZ formalism we actually need not the correlation function, but rather its Fourier transform. E.g. we will need expressions of the form

$$
\begin{align*}
& \left.\int \mathrm{d}^{4} x_{1} e^{i p_{1} \cdot x_{1}} \int \mathrm{~d}^{4} x_{2} e^{i p_{2} \cdot x_{2}} \int \mathrm{~d}^{4} x_{3} e^{-i p_{3} \cdot x_{3}} \int \mathrm{~d}^{4} x_{4} e^{-i p_{4} \cdot x_{4}}\langle 0| \ldots|0\rangle\right|_{\text {Diagram } \times} \\
= & (2 \pi)^{4} \delta^{(4)}\left(p_{1}+p_{2}-q_{3}-q_{4}\right)(-i \lambda) \times  \tag{2.139}\\
& \frac{i}{p_{1}^{2}-m_{0}^{2}+i \epsilon} \frac{i}{p_{2}^{2}-m_{0}^{2}+i \epsilon} \frac{i}{p_{3}^{2}-m_{0}+i \epsilon} \frac{i}{p_{4}^{2}-m_{0}+i \epsilon}
\end{align*}
$$

### 2.8 Disconnected diagrams

A typical diagram contains disconnected pieces, i.e. subdiagrams which are not connected to any of the external points. For example consider the second diagram in Figure 2.4: The double loop in this
case is not connected to any external point and therefore called 'disconnected piece' (or disconnected diagram). A disconnected piece contains only internal points. By contrast, the part of the diagram which is connected to at least one external point is called 'partially connected diagram'.
According to the Feynman rules disconnected pieces appear as overall factors of the Feynman diagram. If we sum up all Feynman diagrams that contribute to

$$
\begin{equation*}
\langle 0| T \Phi\left(x_{1}\right) \ldots \Phi\left(x_{k}\right) e^{-i \frac{\lambda}{N!} \int \Phi^{N} \mathrm{~d}^{4} x}|0\rangle \tag{2.140}
\end{equation*}
$$

the result factorises into the sum of all 'partially connected diagrams' in the above sense multiplied by the sum of all disconnected diagrams. If you don't believe this, just draw a few pages of Feynman diagrams relevant for, say, the 4 -point function in $\phi^{4}$-theory and convince yourself that the infinite sum over all relevant diagrams can be organized this way.
Let $\left\{V_{j}\right\}$ denote the set of all individual disconnected pieces (i.e. $V_{i}$ is a Feynman diagram containing only internal points which itself is connected as a diagram). Then the sum over all diagrams not connected to any of the external points can be organized as a sum over all $V_{i}$, where each $V_{i}$ appears $n_{i}$-times with $n_{i}=0,1, \ldots, \infty$. If some $V_{i}$ appears $n_{i}$-times we must divide its contribution by $n_{i}$ ! to account for the symmetry factor from interchanging identical $V_{i}$. For instance, the sum over all disconnected diagrams contains

$$
1+8+\frac{1}{2!} 88+\frac{1}{3!} 888+\ldots+\text { arbitrary products and number of loops. }
$$

Denoting by $V_{i}$ also the value of the corresponding diagram, we have that the sum over all disconnected diagrams is

$$
\begin{equation*}
\prod_{i} \sum_{n_{i}=0}^{\infty}\left(V_{i}\right)^{n_{i}} \frac{1}{n_{i}!}=\prod_{i} e^{V_{i}}=e^{\Sigma_{i} V_{i}} . \tag{2.141}
\end{equation*}
$$

This is called exponentiation of disconnected pieces. The final result for (2.140) is therefore of the form

$$
\begin{equation*}
\sum(\text { at least partially connected pieces }) \times e^{\Sigma_{i} V_{i}} . \tag{2.142}
\end{equation*}
$$

Now consider the full correlator

$$
\begin{equation*}
\langle\Omega| T \prod_{i} \phi_{i}|\Omega\rangle=\frac{\langle 0| T \prod_{i} \phi_{i} e^{-i \int H_{I} \mathrm{dt}}|0\rangle}{\langle 0| T e^{-i \int H_{I} \mathrm{dt}}|0\rangle} . \tag{2.143}
\end{equation*}
$$

The denominator contains no external points and thus yields precisely

$$
\begin{equation*}
\langle 0| T e^{-i \int H_{\mathrm{I}} \mathrm{~d} t}|0\rangle=e^{\Sigma_{i} V_{i}} \equiv \text { the partition function. } \tag{2.144}
\end{equation*}
$$

This cancels the factor from the nominator and thus

$$
\begin{equation*}
\langle\Omega| T \prod_{i}^{n} \phi_{i}|\Omega\rangle=\sum \text { (all partially connected diagrams with } n \text { external points). } \tag{2.145}
\end{equation*}
$$

Note that the partition function plays an important role also in Statistical Mechanics. In computing the partition function in that context you will encounter many times the same reasoning that organizes the sum over all its contributions in a form similar to the one above. In fact, when discussing the path integral approach to Quantum Field Theory in QFT II we will see that this is no coincidence.

### 2.8.1 Vacuum bubbles

The disconnected diagrams contributing to the partition function are called vacuum bubbles and they have a remarkable physical interpretation: From eq. (2.103) we recall that

$$
\begin{equation*}
e^{\Sigma_{i} V_{i}}=\langle 0| T e^{-i \int \mathrm{~d} t H_{I}}|0\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)}|\langle\Omega \mid 0\rangle|^{2} e^{-i E_{\Omega} 2 T} \tag{2.146}
\end{equation*}
$$

so

$$
\begin{equation*}
-\left.i E_{\Omega} 2 T\right|_{T \rightarrow \infty(1-i \epsilon)}=\sum V_{i}-\underbrace{\log \left(|\langle\Omega \mid 0\rangle|^{2}\right)}_{\text {finite number }} \tag{2.147}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
E_{\Omega}=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{i \sum V_{i}}{2 T} \tag{2.148}
\end{equation*}
$$

because the finite term plays no role in the limit we are taking. $E_{\Omega}$ is the vacuum energy of the interacting theory in the scheme where $E_{0}$ - the vacuum energy of the free theory - has been set to 0 by renormalisation of the original Lagrangian after adding a term $-\frac{E_{0}}{\mathrm{Vol}_{\mathrm{R}^{3}}}$ to $\mathcal{L}$. In terms of the corresponding energy density this means

$$
\begin{equation*}
\frac{E_{\Omega}}{\operatorname{Vol}_{\mathbb{R}^{3}}}=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{i \sum_{j} V_{j}}{2 T \operatorname{Vol}_{\mathbb{R}^{3}}}=i \frac{\sum_{j} V_{j}}{(2 \pi)^{4} \delta^{(4)}(0)} . \tag{2.149}
\end{equation*}
$$

Note that the factor of $(2 \pi)^{4} \delta^{(4)}(0)$ in the denominator cancels a corresponding factor appearing in the computation of each of the $V_{i}$ from an uncancelled overall $\int d^{4} z$ - see footnote 8. This factor represents the familiar IR divergence of the vacuum energy from integration over spacetime. As we will see later when actually computing loops of this type, $\frac{E_{\Omega}}{V_{\mathrm{R}^{3}}{ }^{3}}$ is also UV divergent because of divergent momentum integration for the vacuum loops. This is no surprise as already $\frac{E_{0}}{\mathrm{Vol}_{\mathbb{R}^{3}}}$ was divergent. We organise these divergences order by order in perturbation theory and can renormalise them away by counter-terms, i.e. by adding to the classical Lagrangian terms of form ${ }^{10}$

$$
\begin{equation*}
\mathcal{L} \rightarrow \mathcal{L}-\lambda^{1} V_{0}^{(1)}-\lambda^{2} V_{0}^{(2)}-\ldots, \tag{2.150}
\end{equation*}
$$

just like we did for $E_{0}$ (i.e. at order $\lambda^{0}$ ) in the free theory.

### 2.9 1-particle-irreducible diagrams

As a first application of our formula for the computation of $n$-point functions in an interacting theory we consider perturbative corrections to the propagator in $\Phi^{4}$-theory,

$$
\begin{equation*}
D_{F}(x-y)=\langle\Omega| T \Phi(x) \Phi(y)|\Omega\rangle=\sum_{i} O\left(\lambda^{i}\right) \tag{2.151}
\end{equation*}
$$

Some of the first few Feynman diagrams are shown in Figure 2.5.

[^13]

Figure 2.5: Perturbative corrections.

To organise the sum over Feynman diagrams it is useful to define the concept of 1 - particle-irreducible (1PI) Feynman diagrams. These are diagrams out of which one cannot produce 2 separate non-trivial diagrams (diagrams containing more than just one line) by cutting a single line. E.g. in Figure 2.5 only the 4 th diagram is not 1 PI (one can cut it in 2 by cutting the line between the two loops). It is standard to introduce the notation

$$
\begin{equation*}
\text { (1PI) }=\sum(\text { all non-trivial 1PI diagrams }), \tag{2.152}
\end{equation*}
$$

where it is understood that we do not attach external points to both ends (i.e. no factors $e^{-i p \cdot x}$ or $i /\left(p^{2}-m_{0}^{2}+i \epsilon\right)$ from the left or right of (1PI)
We now define $-i M^{2}\left(p^{2}\right)$ to be the value of (1PI), where $p^{2}$ denotes the in- and outgoing momentum. Clearly computing $-i M^{2}\left(p^{2}\right)$ can be a hard task, and the result will be a complicated function of $p^{2}$. In particular this quantity may be divergent (again due to integration over internal momenta). We will compute $-i M^{2}\left(p^{2}\right)$ explicitly for an electron in Quantum Electrodynamics, in which case it is called self-energy, later in the course and find ways to deal with its divergence systematically. Irrespective of the outcome for $-i M^{2}\left(p^{2}\right)$, the sum over all diagrams contributing to $\langle\Omega| T \Phi(x) \Phi(y)|\Omega\rangle$ can be organised as a geometric series in 1 PI diagrams. Let us compute the Fourier transform $D_{F}\left(p^{2}\right)$ defined such that

$$
\begin{equation*}
D_{F}(x-y)=\int \frac{d^{4} p}{(2 \pi)^{4}} e^{-i p(x-y)} D_{F}\left(p^{2}\right) \tag{2.153}
\end{equation*}
$$

$D_{F}\left(p^{2}\right)$ can be expanded as

$$
\begin{equation*}
D_{F}\left(p^{2}\right)=\frac{i}{\frac{i}{p^{2}-m_{0}^{2}+i \epsilon}}+\frac{i}{\frac{i}{p^{2}-m_{0}^{2}+i \epsilon}} \frac{1}{p^{2}-n_{0}^{2}+i \epsilon}+\cdots \text { (1PI)-1PI)- }+\ldots \tag{2.154}
\end{equation*}
$$

The fact that we are not drawing points at the two ends of diagrams symbolizes that we are omitting the factors of $e^{i p x}$ or $e^{-i p y}$ as these are, by definition, not part of $D_{F}\left(p^{2}\right)$. The Fourier transformed
propagator is then

$$
\begin{align*}
D_{F}\left(p^{2}\right) & =\frac{i}{p^{2}-m_{0}^{2}+i \epsilon}+\frac{i}{p^{2}-m_{0}^{2}+i \epsilon}\left(-i M^{2}\left(p^{2}\right)\right) \frac{i}{p^{2}-m_{0}^{2}+i \epsilon}+\ldots \\
& =\frac{i}{p^{2}-m_{0}^{2}+i \epsilon}[1+\underbrace{\frac{M^{2}\left(p^{2}\right)}{p^{2}-m_{0}^{2}+i \epsilon}}_{\equiv q}+\left(\frac{M^{2}\left(p^{2}\right)}{p^{2}-m_{0}^{2}+i \epsilon}\right)^{2}+\ldots]  \tag{2.155}\\
& =\frac{i}{p^{2}-m_{0}^{2}+i \epsilon} \frac{1}{1-\frac{M^{2}\left(p^{2}\right)}{p^{2}-m_{0}^{2}+i \epsilon}} .
\end{align*}
$$

Here we assumed that $q<1$ so that the geometric series $\sum_{n=0}^{\infty} q^{n}$ converges to $\frac{1}{1-q}$. Once again, we will deal with the explicit form of $-i M^{2}\left(p^{2}\right)$ later in QED and justify this assumption. This procedure is called Dyson resummation and yields the resummed propagator

$$
\begin{equation*}
D_{F}\left(p^{2}\right)=\frac{i}{p^{2}-\left(m_{0}^{2}+M^{2}\left(p^{2}\right)\right)+i \epsilon} \equiv \square \tag{2.156}
\end{equation*}
$$

Now, $-i M^{2}\left(p^{2}\right)$ is a complicated function of $p^{2}$ that we compute order by order in $\lambda$. We can then extract the first analytic pole in $D_{F}\left(p^{2}\right)$ and call it $m^{2}$. Then

$$
\begin{equation*}
D_{F}\left(p^{2}\right)=\frac{i Z}{p^{2}-m^{2}+i \epsilon}+\text { terms regular at } m^{2} \tag{2.157}
\end{equation*}
$$

for some $Z$. I.e. $Z$ is by definition the residue (in the above sense) of $D_{F}\left(p^{2}\right)$ at its first analytical pole $m^{2}$. This way we can compute

- $m^{2}$ - the physical 1-particle-mass $\equiv$ pole mass,
- $Z$ - the wavefunction renormalisation,
perturbatively to given order in $\lambda$. This is a beautiful result because we can now understand quantitatively why the mass of the 1-particle momentum eigenstates in an interacting theory differs from $m_{0}$ : The reason are the self-interactions of the field, which are resummed as above to shift the pole of the full propagator from $m_{0}$ to $m$. In particular this picture justifies our assertion made in the context of asymptotic in- and out-states that the asymptotic states behave as free particles, but with fully 'renormalised mass' $m \neq m_{0}$ : By sending the particles infinitely far apart from each other we effectively 'switch off' the interactions between the different particles, but we cannot switch off the interactions of the asymptotic particles with themselves (or rather of the field with itself). These are precisely the 1PI contributions to $D_{F}\left(p^{2}\right)$ and thus the in-and out-states do have the fully resummed mass $m^{2} \neq m_{0}^{2}$.


### 2.10 Scattering amplitudes

Recall the LSZ-reduction formula for $r$ in- and $n$ out-states

$$
\begin{align*}
& \prod_{k=1}^{n} \int \mathrm{~d}^{4} y_{k} e^{i p_{k} \cdot y_{k}} \prod_{l=1}^{r} \int \mathrm{~d}^{4} x_{l} e^{-i q_{l} \cdot x_{l}}\langle\Omega| T \Phi\left(y_{1}\right) \ldots \Phi\left(x_{1}\right) \ldots|\Omega\rangle \\
& =\left.\prod_{k} \frac{i Z^{1 / 2}}{p_{k}^{2}-m^{2}+i \epsilon} \prod_{l} \frac{i Z^{1 / 2}}{q_{l}^{2}-m^{2}+i \epsilon}\left\langle p_{1} \ldots\right| S\left|q_{1} \ldots\right\rangle\right|_{\text {connected }} \tag{2.158}
\end{align*}
$$

where all $q_{l}$ and $p_{k}$ are on-shell. Now, recall that $\left.\left\langle p_{1} \ldots\right| S\left|q_{1} \ldots\right\rangle\right|_{\text {connected }}$ itself cannot be proportional to $p_{k}^{2}-m^{2}$ or $q_{l}^{2}-m^{2}$ because otherwise it is zero on-shell and thus no scattering occurs; likewise $\left.\left\langle p_{1} \ldots\right| S\left|q_{1} \ldots\right\rangle\right|_{\text {connected }}$ cannot contain any factors of $\left(p_{k}^{2}-m^{2}\right)^{-1}$ or $\left(q_{l}^{2}-m^{2}\right)^{-1}$ because then it would be divergent on-shell, in contradiction with its definition as a quantum mechanical amplitude. As a result only those Feynman diagrams are relevant with exactly $(n+r)$ poles at $m^{2}$ in the above sense. It is not hard to see that these are precisely the ones that contribute to the fully connected correlation function. A fully connected correlation function has the structure displayed in Figure 2.6.


Figure 2.6: Fully connected correlation function with $r$ in and $n$ out states.

Here $\times$ denotes the fully resummed propagator from lines connected to external points. The fully resummed propagator would appear if we computed to all orders in perturbation theory. More realistically, to a given finite order in perturbation theory we should think of the propagator as computed to suitable order.

The big blop in the middle denotes the amputated correlator. By amputated correlator we mean the Feynman diagram after cutting off all external legs carrying $\times$ Since each external leg carries a factor of $\frac{i Z}{p^{2}-m^{2}+i \epsilon}$ near $m^{2}$ if $p^{2}$ is on-shell, all $(n+r)$ external legs yield together yield the right singularity structure.
By contrast, partially connected diagrams such as the one in Figure 2.7 carry fewer factors of $\frac{i}{p^{2}-m^{2}+i \epsilon}$ and thus do not contribute to $\left.\left\langle p_{1} \ldots\right| S\left|q_{1} \ldots\right\rangle\right|_{\text {connected }}$. Therefore the final result for the computation

## of scattering amplitudes is

$$
\begin{align*}
& \left.\left\langle p_{1}, \ldots, p_{n}\right| S\left|q_{1}, \ldots, q_{r}\right\rangle\right|_{\text {connected }}=(\sqrt{Z})^{n+r} \times  \tag{2.159}\\
& \left.\left(\left.\prod_{k} \int \mathrm{~d}^{4} y_{k} e^{i p_{k} \cdot y_{k}} \prod_{l} \int \mathrm{~d}^{4} x_{l} e^{-i q_{l} \cdot x_{l}}\langle\Omega| T \Phi\left(y_{1}\right) \ldots \Phi\left(x_{1}\right) \ldots|\Omega\rangle\right|_{\text {fully connected }}\right)\right|_{\text {Amputated }}
\end{align*}
$$

Note that the wavefunction renormalisation factor $Z$ itself is of the form $1+O(\lambda)$ in perturbation theory, so to leading order in the coupling constant $\lambda$ the $Z$-factors play no role as only $O(\lambda)$ diagrams can be fully connected.


Figure 2.7: Partially connected diagramm, i.e. $(n-2)$ and $(r-2)$.

As an example we consider again the $\Phi^{4}$-theory to leading order. The only fully connected diagram at $O(\lambda)$ is the first diagram in Figure 2.4. We make a consistent choice of ingoing versus outgoing momenta and recall the result

$$
\begin{align*}
& \left.\langle\Omega| T \Phi\left(x_{1}\right) \Phi\left(x_{2}\right) \Phi\left(y_{1}\right) \Phi\left(y_{2}\right)|\Omega\rangle\right|_{\times}= \\
& =\prod_{j=1}^{4} \int \frac{\mathrm{~d}^{4} k_{j}}{(2 \pi)^{4}} e^{-i k_{1} \cdot x_{1}} e^{-i k_{2} \cdot x_{2}} e^{i k_{3} \cdot y_{1}} e^{i k_{4} \cdot y_{2}}  \tag{2.160}\\
& \quad \times(-i \lambda)(2 \pi)^{4} \delta^{(4)}\left(k_{1}+k_{2}-k_{3}-k_{4}\right) \prod_{l=1}^{4} \frac{i}{k_{l}^{2}-m_{0}^{2}+i \epsilon}
\end{align*}
$$

To compute $\left.\left\langle p_{1}, p_{2}\right| S\left|q_{1}, q_{2}\right\rangle\right|_{\text {connected }}$ we take first the Fourier transform as

$$
\begin{align*}
& \int \mathrm{d}^{4} x_{1} e^{-i q_{1} \cdot x_{1}} \int \mathrm{~d}^{4} x_{2} e^{-i q_{2} \cdot x_{2}} \int \mathrm{~d}^{4} y_{1} e^{i p_{1} \cdot y_{1}} \int \mathrm{~d}^{4} y_{2} e^{i p_{2} \cdot y_{2}} \\
& \quad \times\left.\langle\Omega| T \Phi\left(x_{1}\right) \Phi\left(x_{2}\right) \Phi\left(y_{1}\right) \Phi\left(y_{2}\right)|\Omega\rangle\right|_{\times}  \tag{2.161}\\
& =(-i \lambda)(2 \pi)^{4} \delta^{(4)}\left(q_{1}+q_{2}-p_{1}-p_{2}\right) \prod_{j=1}^{2} \frac{i}{q_{j}^{2}-m_{0}^{2}+i \epsilon} \frac{i}{p_{j}^{2}-m_{0}^{2}+i \epsilon}
\end{align*}
$$

and then we amputate by discarding all propagators from external lines. The result is

$$
\begin{equation*}
\left.\left\langle p_{1}, p_{2}\right| S\left|q_{1}, q_{2}\right\rangle\right|_{\text {connected, } O(\lambda)}=(-i \lambda)(2 \pi)^{4} \delta^{(4)}\left(q_{1}+q_{2}-p_{1}-p_{2}\right), \tag{2.162}
\end{equation*}
$$

where we have discarded the factor $\sqrt{Z}^{4}$ because, as discussed, it does not contribute to the leading order result.

### 2.10.1 Feynman-rules for the $S$-matrix

This procedure can be summarised in general in terms of the Feynman-rules for the computation of $\left.\left\langle p_{1}, \ldots, p_{n}\right| S\left|q_{1}, \ldots, q_{r}\right\rangle\right|_{\text {connected }}$ as follows:

- Draw the relevant fully connected Feynman diagrams with $(n+r)$ external points to given order in $\lambda$.
- Assign ingoing momenta $q_{l}$ and outgoing momenta $p_{k}$ and label momenta of internal lines with $k_{j}$.
- Each vertex carries $(-i \lambda)(2 \pi)^{4} \delta^{(4)}$ ( $\sum$ ingoing momenta - outgoing momenta).
- Each internal line carries $\frac{i}{k_{j}^{2}-m_{0}^{2}+i \epsilon}$.
- Integrate over all internal momenta $\prod_{j} \int \frac{\mathrm{~d}^{4} k_{j}}{(2 \pi)^{4}}$ and divide by the symmetry factor.
- Sum up all diagrams and multiply by $(\sqrt{Z})^{n+r}$ to given order in $\lambda$.


## Deeper interpretation of the Feynman rules

The connection between correlation functions and $S$-matrix elements gives Feynman diagrams an intuitive physical meaning: In Figure 2.8 two particles come in and interact at $Z_{1}$ to form two socalled "virtual" particles, which in turn join again at $Z_{2}$ and form two outgoing states. The obvious


Figure 2.8: Feynman graph with two virtual particles.
interpretation of the Feynman diagram elements is therefore:

- A line - corresponds to the worldline of a particle.
$\bullet \longrightarrow e^{-i p \cdot x}$ is the wavefunction for a momentum-eigenstate.
- A vertex $Z$ is a localised interaction at a spacetime point $Z$.
- Summing up diagrams and integrating over $\int \mathrm{d}^{4} z$ amounts to coherently summing up the QM probability amplitudes for all possible processes - called "channels" - with the same macroscopic result. In particular, this will lead to QM interference between the different channels.

Intermediate particles, e.g. those running in the loop as $k_{1}$ and $k_{2}$ in Figure 2.8, are called virtual because they are generally off-shell: This means that these particles do in general not satisfy the relation $k_{i}^{2}-m_{0}^{2}=0$. E.g. for the loop-diagram above we integrate over all internal momenta (where the delta-functions at the vertices implies that some of these integrals become trivial). In the present diagram momentum conservation at the first vertex implies that

$$
\begin{equation*}
q_{1}+q_{2}=k_{1}+k_{2} \Leftrightarrow k_{2}=q_{1}+q_{2}-k_{1}, \tag{2.163}
\end{equation*}
$$

where $k_{1}$ is free and thus integrated over. For most values of $k_{1}$ we have $k_{1}^{2}-m_{0}^{2} \neq 0$. We can think of this as a quantum-mechanical violation of the energy-momentum relation because $E_{j}^{2} \neq$ $\vec{k}_{j}^{2}+m_{0}^{2}$ for virtual particles. In Quantum Mechanics this is allowed for sufficiently short times - in fact this precisely what is meant by the energy-time uncertainty relation arising in quantum mechanical perturbation theory. Note, however, that 4-momentum conservation is guaranteed at each vertex due to the factors

$$
\begin{equation*}
(2 \pi)^{4}\left(\sum \text { ingoing momenta }- \text { outgoing momenta }\right) . \tag{2.164}
\end{equation*}
$$

## LSZ versus Interaction Picture Perturbation Theory

On assignment 7 will work out a detailed a comparison of the LSZ approach to scattering as completed in this chapter and a more naive approach based on quantum mechanical Interaction Picture perturbation theory.

### 2.11 Cross-sections

$S$-matrix elements are in general of the form

$$
\begin{equation*}
\langle f| S|i\rangle=\underbrace{\delta_{f i}}_{\text {no scattering }}+\underbrace{i(2 \pi)^{4} \delta^{(4)}\left(p_{f}-p_{i}\right)}_{\text {momentum conservation }} \cdot \underbrace{\mathcal{M}_{f i}}_{=: \text {scattering amplitude }} \tag{2.165}
\end{equation*}
$$

In particular, momentum conservation is a consequence of the insertion of delta-functions at each vertex in the Feynman diagrams. The first summand $\delta_{f i}$ accounts for the situation that the initial and final states are identical, in which case no scattering occurs. In the sequel we will usually exclude this possibility and focus on non-trivial scattering events.
The QM probability for scattering of a given initial state $\{|i\rangle\}$ into a range of final states $\{|f\rangle\}$ is

$$
\begin{equation*}
\left.\mathcal{P}_{|i\rangle \rightarrow\{|f\rangle\rangle}=\sum_{|f\rangle \in \||f\rangle\rangle}|\langle f| S| i\right\rangle\left.\right|^{2} . \tag{2.166}
\end{equation*}
$$

Assume now that $|i\rangle \notin\{|f\rangle\}$ (so that in particular the term $\delta_{f i}$ is irrelevant), then

$$
\begin{equation*}
\mathcal{P}_{|i\rangle \rightarrow| | f\rangle\}}=\sum_{|f\rangle \in| | f\rangle\}} \underbrace{\left[\left(p^{(4)}\right) p_{i}\right]^{2}}_{=(2 \pi)^{4} \delta^{(4)}\left(p_{f}-p_{i}\right) \underbrace{\left[(2 \pi)^{4} \delta^{4} \delta^{(4)}(0)\right.}_{=V_{\mathbb{R}}, 3}}\left|\mathcal{M}_{f i}\right|^{2} . \tag{2.167}
\end{equation*}
$$

We define the transition rate as the probability normalised per spacetime volume,

$$
\begin{equation*}
\omega_{f i}=\frac{\mathcal{P}_{|i\rangle \rightarrow| | f\rangle\}}}{\text { unit time } \times \text { unit volume }} \tag{2.168}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\omega_{f i}=\sum_{|f\rangle \in \backslash|f\rangle\rangle}(2 \pi)^{4} \delta^{(4)}\left(p_{f}-p_{i}\right)\left|\mathcal{M}_{f i}\right|^{2} \tag{2.169}
\end{equation*}
$$

For scattering into $N$ identical particles this is

$$
\begin{equation*}
\omega_{f i}=\frac{1}{N!} \prod_{n=1}^{N} \int \frac{\mathrm{~d}^{3} k_{n}}{(2 \pi)^{3}} \frac{1}{2 E_{n}}(2 \pi)^{4} \delta^{(4)}\left(\sum_{i} p_{i}-\sum_{n} k_{n}\right)\left|\mathcal{M}_{f i}\right|^{2}, \tag{2.170}
\end{equation*}
$$

where the symmetry factor $\frac{1}{N!}$ accounts for the indistinguishability of the $N$ identical particles. Now, a typical scattering experiment is of the form

$$
\text { Incoming beam of density } \rho_{B} \text { and length } l_{B} \rightarrow \text { Target } A \text { of density } \rho_{A} \text { and length } l_{A},
$$

where $\rho_{A, B}$ denote the respective number densities (i.e. number of particles per volume). The number (\#) of scattered particles is proportional ${ }^{11}$ to $l_{A} l_{B} \int \mathrm{~d}^{2} x \rho_{a}(x) \rho_{B}(x)$. The factor of proportionality has dimension [Area] and is called cross-section $\sigma$,

$$
\begin{equation*}
\# \text { of events }=: \sigma l_{A} l_{B} \int \mathrm{~d}^{2} x \rho_{A}(x) \rho_{B}(x) . \tag{2.171}
\end{equation*}
$$

If $\rho_{A}$ and $\rho_{B}$ are constant and both beams overlap over the area $\mathcal{A}$, then

$$
\begin{equation*}
\sigma=\frac{\# \text { of scattered particles }}{l_{A} l_{B} \rho_{A} \rho_{B} \mathcal{A}} \tag{2.172}
\end{equation*}
$$

This suggests the following intuitively clear interpretation of the meaning of $\sigma$ : The cross-section $\sigma$ is the effective area of the beam $B$ that participates in the scattering.

Consider now a $2 \rightarrow N$ scattering process such that all out-going states $\left|k_{j}\right\rangle$ are momentum eigenstates. The initial states are the momentum eigenstates $\left|p_{A}\right\rangle$ and $\left|p_{B}\right\rangle$. Let us go to the rest frame of

[^14]the target $A$ called lab frame. Suppose for simpilcity that the beam $B$ hits the target $A$ over its entire area $\mathcal{A}$ and that the densities of beam and target do not vary much. Then
\[

$$
\begin{equation*}
\frac{\# \text { of events }}{\text { Volume } \times \text { Time }}=\frac{\sigma \rho_{A} \rho_{B} l_{A} l_{B} \mathcal{A}}{\text { Volume } \times \text { Time }}=\sigma \rho_{A} \rho_{B}\left|\vec{v}_{B}^{L}\right|, \tag{2.173}
\end{equation*}
$$

\]

where the relevant interaction volume is is just Volume $=\mathcal{A} l_{A}$ and the relevant interaction time is what it takes for the beam to pass over a given slice in the target orthogonal to the beam direction, i.e. $\frac{l_{B}}{\text { Time }}=\left|\vec{v}_{B}^{(L)}\right|$. Sometimes it is also useful to think in terms of the beam flux $\mathcal{F}$ defined as

$$
\begin{equation*}
\rho_{B}\left|\vec{v}_{B}^{L}\right|=: \mathcal{F} . \tag{2.174}
\end{equation*}
$$

Since $E=\gamma m$ and $\vec{p}=m \gamma \vec{v}$ we can use

$$
\begin{equation*}
\left|\vec{v}_{B}^{L}\right|=\frac{\left|\vec{p}_{B}^{L}\right|}{E_{B}} \tag{2.175}
\end{equation*}
$$

to obtain

$$
\begin{equation*}
\frac{\# \text { events }}{\text { Volume } \times \text { Time }}=\sigma \rho_{A} \rho_{B} \frac{\left|\vec{p}_{B}^{L}\right|}{E_{B}}=\sigma \rho_{A} \mathcal{F} \tag{2.176}
\end{equation*}
$$

What we actually compute in evaluating the S-matrix elements are not absolute numbers of events, but rather QM probabilites. In order make the transition between both these concepts we replace the particle number densities $\rho_{A}$ and $\rho_{B}$ by the quantum mechanical densities of the corresponding single particles states, i.e. the probabilities of finding a particle $A$ or $B$ per given volume. If $\left\langle p_{A} \mid p_{A}\right\rangle=1$ this probability would be $1 / V_{\mathbb{R}^{3}}=\left\langle p_{A} \mid p_{A}\right\rangle / V_{\mathbb{R}^{3}}$. With our normalisation of momentum states

$$
\begin{equation*}
\left\langle p_{A} \mid p_{A}\right\rangle=(2 \pi)^{3} 2 E_{A} \delta^{(3)}\left(\vec{p}_{A}-\vec{p}_{A}\right)=2 E_{A} V_{\mathbb{R}_{3}} \tag{2.177}
\end{equation*}
$$

we must replace

$$
\begin{equation*}
\rho_{A} \rightarrow \frac{\left\langle p_{A} \mid p_{A}\right\rangle}{V_{\mathbb{R}^{3}}}=2 E_{A} \quad \text { and } \quad \rho_{b} \rightarrow 2 E_{B} \tag{2.178}
\end{equation*}
$$

Since we are in the rest frame of $A$ we have $2 E_{A}=2 m$. Therefore in the lab frame $\left(\vec{v}_{A}^{(L)}=0\right)$

$$
\begin{equation*}
\omega_{f i}=\frac{\mathrm{QM} \text { probability }}{\text { Volume } \times \text { Time }}=4 \sigma m p_{B}^{(L)} \tag{2.179}
\end{equation*}
$$

Note that

$$
\begin{equation*}
4 m p_{B}^{(L)}=4 E_{A} E_{B}\left|\vec{v}_{B}^{(L)}\right|=4 E_{A} E_{B}\left|\vec{v}_{A}^{(L)}-\vec{v}_{B}^{(L)}\right| \tag{2.180}
\end{equation*}
$$

in the lab frame. Finally, one can show that $4 E_{A} E_{B}\left|\vec{v}_{A}^{(L)}-\vec{v}_{B}^{(L)}\right|$ is invariant under Lorentz boosts in direction $A \rightarrow B$. Therefore $\omega_{f i}=4 E_{A} E_{B}\left|\vec{v}_{A}-\vec{v}_{B}\right| \sigma$ is the correct general expression for the transition rate valid in any frame ${ }^{12}$. To conclude, the differential cross-section is

$$
\begin{equation*}
\mathrm{d} \sigma=\frac{(2 \pi)^{4}}{4 E_{A} E_{B}\left|\vec{v}_{A}-\vec{v}_{B}\right|} \mathrm{d} \Pi_{N} \delta^{(4)}\left(p_{a}+p_{B}-\sum_{i} k_{i}\right)\left|\mathcal{M}_{f i}\right|^{2} \tag{2.181}
\end{equation*}
$$

[^15]where
\[

$$
\begin{equation*}
\mathrm{d} \Pi_{N} \equiv \frac{1}{N!} \prod_{n=1}^{N} \int \frac{\mathrm{~d}^{3} k_{n}}{(2 \pi)^{3}} \frac{1}{2 E_{n}} . \tag{2.182}
\end{equation*}
$$

\]

For example consider $2-2$ scattering as shown in Figure 2.9. On Assignment 7 we show that if all


Figure 2.9: $2-2$ scattering.
four particles have the same mass, the differential cross-section is

$$
\begin{equation*}
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega_{3}}=\frac{1}{2!} \frac{1}{64 \pi^{2}} \frac{1}{s}|\mathcal{M}|^{2}, \tag{2.183}
\end{equation*}
$$

where the relativistically invariant Mandelstam variable

$$
\begin{equation*}
s=\left(p_{1}+p_{2}\right)^{2}=\left(p_{3}+p_{4}\right)^{2} \tag{2.184}
\end{equation*}
$$

appears. Note that in $\Phi^{4}$-theory $|\mathcal{M}|^{2}=\lambda^{2}$ to first order in $\lambda$, i.e. for the maximally localised, pointlike interaction given by the first Feynman diagram in Figure 2.4. Therefore

$$
\begin{equation*}
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega} \sim \frac{1}{s} \tag{2.185}
\end{equation*}
$$

to first order in $\lambda$. This is a famous result, which in fact holds more generally:

The differential cross-section for hard scattering off a pointlike target (i.e. a target with no substructure of length $l \geq 1 / \sqrt{s}$ ) falls off as $1 / s$.

This characteristic behaviour as observed in deep inelastic scattering experiments with hardons was a crucial clue to the the parton structure of hadrons as you will surely recall from your particle physics course.

## Chapter 3

## Quantising spin $\frac{1}{2}$-fields

### 3.1 The Lorentz algebra so $(1,3)$

## Relativistic fields as representations of the Lorentz algebra

Relativistic fields are classified by their behaviour under Lorentz transformations

$$
\begin{equation*}
x^{\mu} \mapsto x^{\prime \mu}=\Lambda_{v}^{\mu} x^{v}, \quad \Lambda_{v}^{\mu} \in S O(1,3) \tag{3.1}
\end{equation*}
$$

Quite generally a field $\Phi^{a}(x)$ transforms as a representation ${ }^{1}$ of the Lorentz group $S O(1,3)$ - or, more precisely, as we will see momentarily, of its double cover $\operatorname{Spin}(1,3)$. This means that the classical field $\Phi^{a}(x)$ is a map

$$
\begin{align*}
\Phi^{a}: \mathbb{R}^{1,3} & \rightarrow V,  \tag{3.2}\\
x & \mapsto \Phi^{a}(x), \quad a=1, \ldots, \operatorname{dim}(V) \tag{3.3}
\end{align*}
$$

with $V$ a vector space such that the specific transformation behaviour of the field $\Phi^{a}(x)$ under (3.1) is given by

$$
\begin{equation*}
\Phi^{a}(x) \mapsto R(\Lambda)^{a}{ }_{b} \Phi^{b}\left(\Lambda^{-1} x^{\prime}\right)=R(\Lambda)^{a}{ }_{b} \Phi^{b}(x) \tag{3.4}
\end{equation*}
$$

Here for every element $\Lambda \in S O(1,3)$ the object $R(\Lambda)$ is an automorphism (an invertible linear map) acting on $V$ in a manner compatible with the group action of $S O(1,3)$. Specifying the representation in which the field transforms amounts to assigning such an automorphism $R$ to every element $\Lambda \in$ $S O(1,3)$, where compatibility with the group action of $S O(1,3)$ means that $^{2}$

$$
\begin{align*}
R\left(\Lambda_{2}\right)^{a}{ }_{b} R\left(\Lambda_{1}\right)^{b}{ }_{c} & =R\left(\Lambda_{2} \Lambda_{1}\right)_{c}^{a} \\
R\left(\Lambda^{-1}\right)_{b}^{a} & =\left(R(\Lambda)^{-1}\right)_{b}^{a} \tag{3.5}
\end{align*}
$$

[^16]The vector space $V$ is also called representation space. Note that the index $a$ refers to the components of the fields with respect to a basis of the vector space $V$ - called the representation space -, while $\mu, \nu$ are spacetime indices. These two are in general completely different objects.

Consider the following examples:

- For a real (or complex) scalar field $\Phi(x)$ the representation space $V$ is just $\mathbb{R}$ (or $\mathbb{C}$ ) and

$$
\begin{equation*}
R(\Lambda)=\mathbb{1} \forall \Lambda . \tag{3.6}
\end{equation*}
$$

The scalar field is said to transform in the trivial or scalar representation and describes particles with spin 0 .

- A vector field $A^{\mu}(x)$ (e.g. the gauge potential of electro-magnetism) transforms in the vector representation: The representation space $V$ is identified with spacetime $\mathbb{R}^{1,3}$ itself (or rather its tangent space) and

$$
\begin{equation*}
R(\Lambda)_{v}^{\mu}=\Lambda_{v}^{\mu} \forall \Lambda . \tag{3.7}
\end{equation*}
$$

Since here $V=\mathbb{R}^{1,3}$, we have in this case that $a, b=\mu, \nu$. A field in the vector representation describes particles with spin 1 (as we will see later in this course when quantising such vector fields).

An important concept in representation theory is that of an irreducible representation (irrep): An Irreducible representation is one whose representation space does not split into a direct sum of two vector spaces in a manner respected by the group action, i.e. it is not possible to find a basis of the representation space in which

$$
R(\Lambda)=\left(\begin{array}{c|c}
\star & 0  \tag{3.8}\\
\hline 0 & \star
\end{array}\right) \forall \Lambda .
$$

Irreducible representations form the building blocks of which larger representations can be formed by considering direct sums of representation spaces.

Spin $1 / 2$ particles are described by fields in the spinor representation. To find this representation our starting point will not be the irreducible representations of the Lorentz group $S O(1,3)$, but of

The Lorentz algebra so (1,3)
This is the algebra of infinitesimal Lorentz transformations connected to the identity. This is analogous to the precedure applied in quantum mechanics to find spin $\frac{1}{2}$ representations of the algebra of spatial rotations $s o(3) \simeq s u(2)$.

A Lorentz transformation can be written infinitesimally as

$$
\begin{equation*}
\Lambda_{v}^{\mu}=\delta_{v}^{\mu}+\omega_{v}^{\mu} . \tag{3.9}
\end{equation*}
$$

Since the metric must be invariant under a Lorentz transformation in the sense that

$$
\begin{equation*}
\Lambda_{\alpha}^{\mu} \Lambda_{\beta}^{v} \eta^{\alpha \beta}=\eta^{\mu v} \tag{3.10}
\end{equation*}
$$

the infinitesimal transformations satisfy $\omega^{\mu \nu}=-\omega^{\nu \mu}$. We introduce a basis of antisymmetric $4 \times 4$ matrices into which we would like to expand $\omega^{\mu \nu}$. This basis contains 6 elements. Let the basis be

$$
\begin{equation*}
\left(M^{A}\right)^{\mu \nu} \text { with } A=1, \ldots, 6 \tag{3.11}
\end{equation*}
$$

Then

$$
\begin{equation*}
\omega^{\mu \nu}=-\frac{i}{2} \Omega_{A}\left(M^{A}\right)^{\mu \nu} \tag{3.12}
\end{equation*}
$$

for some $\Omega_{A}$ with summation over $A$ understood. We introduce the double-index notation $A=\rho \sigma$ with $\rho, \sigma=0,1,2,3$ such that

$$
\begin{equation*}
\left(M^{\rho \sigma}\right)^{\mu \nu}=-\left(M^{\sigma \rho}\right)^{\mu \nu} \tag{3.13}
\end{equation*}
$$

is antisymmetric in $\sigma \rho$. Therefore

$$
\begin{equation*}
\omega^{\mu \nu}=-\frac{i}{2} \Omega_{\rho \sigma}\left(M^{\rho \sigma}\right)^{\mu \nu} \tag{3.14}
\end{equation*}
$$

So an infinitesimal Lorentz transformation is

$$
\begin{equation*}
\Lambda_{v}^{\mu}=\delta_{v}^{\mu}-\frac{i}{2} \Omega_{\rho \sigma}\left(M^{\rho \sigma}\right)_{v}^{\mu} \tag{3.15}
\end{equation*}
$$

A finite transformation is given by applying the infnitesimal version $N$ times for $N \rightarrow \infty$ in the sense

$$
\begin{align*}
\Lambda_{v}^{\mu}\left(M^{\rho \sigma}\right) & =\lim _{N \rightarrow \infty}\left(\delta_{v}^{\mu}-\frac{1}{N} \frac{i}{2} \Omega_{\rho \sigma}\left(M^{\rho \sigma}\right)^{\mu}{ }_{v}\right)^{N}  \tag{3.16}\\
& =\left[e^{-\frac{i}{2} \Omega_{\rho \sigma}\left(M^{\rho \sigma}\right)}\right]_{v}^{\mu}
\end{align*}
$$

This is what is meant by saying that the $M^{\rho \sigma}$ are the generators of the Lie group $S O(1,3)$, or equivalently form a basis of the Lie algebra $\operatorname{so}(1,3)$.
The structure of the Lie algebra $s o(1,3)$ is encoded in the commutation relations of its basis elements. To find these one expands the defining relation of the underlying Lie algebra $S O(1,3)$

$$
\begin{equation*}
\Lambda_{\alpha}^{\mu} \Lambda_{\beta}^{v} \eta^{\alpha \beta}=\eta^{\mu v} \tag{3.17}
\end{equation*}
$$

to second order in $\Omega$ to conclude that

$$
\begin{equation*}
\left[M^{\rho \sigma}, M^{\tau v}\right]=-i\left(\eta^{\rho v} M^{\sigma \tau}+\eta^{\sigma \tau} M^{\rho v}-\eta^{\rho \tau} M^{\sigma v}-\eta^{\sigma v} M^{\rho \tau}\right) \tag{3.18}
\end{equation*}
$$

This defines the structure constants of the Lie-algebra $\operatorname{so}(1,3)$. Recall that in general a basis $T^{A}$ of a Lie algebra satisfies

$$
\begin{equation*}
\left[T^{A}, T^{B}\right]=\sum_{C} i f^{A B}{ }_{C} T^{C} \tag{3.19}
\end{equation*}
$$

where $f^{A B}{ }_{C}$ are the structure constants. They obey

- $f^{A B}{ }_{C}=-f^{B A}{ }_{C}$
- and the Jacobi identity

$$
\begin{equation*}
f^{A B}{ }_{D} f^{C D}{ }_{E}+f^{C A}{ }_{D} f^{B D}{ }_{E}+f^{B C}{ }_{D} f^{A D}{ }_{E}=0 . \tag{3.20}
\end{equation*}
$$

We view the commutator (3.18) as the defining relation for abstract objects $M^{\rho \sigma}$. An $n$-dimensional representation of the Lie algebra $s o(1,3)$ is an assignment that associates to each $M^{\rho \sigma}$ an invertible map

$$
\begin{equation*}
\left(R^{\rho \sigma}\right)^{a}{ }_{b}=R\left(M^{\rho \sigma}\right)^{a}{ }_{b} \tag{3.21}
\end{equation*}
$$

acting on an $n$-dimensional vector space $V$ such that the same relation (3.18) holds for the representation matrices $\left(R^{\rho \sigma}\right)_{b}^{a}$. The indices $a, b=1, \ldots, n$ refer to a basis $e_{a}$ of $V$ in the sense that $v \in V$ is expanded as $v=v^{a} e_{a}$. Note again that in general $a$ and $b$ are unrelated to $\mu$ and $v$.
Consider as a special case the vector representation by choosing $V=\mathbb{R}^{1,3}$ (viewed as the tangent space to spacetime). In this case we do identify $a, b \equiv \mu, \nu$ and set

$$
\begin{equation*}
\left(\mathcal{J}^{\rho \sigma}\right)^{\mu \nu}:=i\left(\eta^{\rho \mu} \eta^{\sigma \nu}-\eta^{\rho \nu} \eta^{\sigma \mu}\right) . \tag{3.22}
\end{equation*}
$$

One can check that this indeed satisfies the so( 1,3 ) relations (3.18) and that

$$
\begin{equation*}
\omega^{\mu \nu}=-\frac{i}{2} \omega_{\rho \sigma}\left(\mathcal{J}^{\rho \sigma}\right)^{\mu \nu} . \tag{3.23}
\end{equation*}
$$

Therefore

$$
\begin{equation*}
\Lambda_{\nu}^{\mu}=\left[e^{-\frac{i}{2} \omega_{\rho \sigma} \mathcal{T}^{\rho \sigma}}\right]_{v}^{\mu} \tag{3.24}
\end{equation*}
$$

We can use this result to deduce the matrix representation of spatial rotations by an angle $\alpha$ around an axis $\vec{n}$. The corresponding $\omega_{\mu \nu}$ can be written as

$$
\begin{equation*}
\omega_{i j}=\alpha \epsilon_{i j k} n^{k} . \tag{3.25}
\end{equation*}
$$

E.g. for a rotation around the $x^{1}$-axis we have

$$
\vec{n}=\left(\begin{array}{l}
1  \tag{3.26}\\
0 \\
0
\end{array}\right) \Rightarrow \omega_{\rho \sigma}=\left(\begin{array}{cccc}
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & \alpha \\
0 & 0 & -\alpha & 0
\end{array}\right)
$$

The corresponding infinitesimal rotations are then of the form

$$
\Lambda_{v}^{\mu}=\delta^{\mu}{ }_{v}+\omega^{\mu}{ }_{v}=\left(\begin{array}{cccc}
1 & 0 & 0 & 0  \tag{3.27}\\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & -\alpha \\
0 & 0 & \alpha & 1
\end{array}\right),
$$

whose finite version is

$$
\Lambda_{v}^{\mu}=\left(\begin{array}{cccc}
1 & 0 & 0 & 0  \tag{3.28}\\
0 & 1 & 0 & 0 \\
0 & 0 & \cos \alpha & -\sin \alpha \\
0 & 0 & \sin \alpha & \cos \alpha
\end{array}\right)
$$

### 3.2 The Dirac spinor representation

To find the spinor representation of $\operatorname{so}(1,3)$ we start from the Clifford algebra Cliff( 1,3 ) defined as the algebra spanned by $n \times n$-matrices $\left(\gamma^{\mu}\right)^{A}{ }_{B}, \mu=0,1,2,3$ and $A, B=1, \ldots, n$ such that the anti-commutator is

$$
\begin{equation*}
\left\{\gamma^{\mu}, \gamma^{\nu}\right\}:=\gamma^{\mu} \gamma^{\nu}+\gamma^{\nu} \gamma^{\mu}=2 \eta^{\mu v} \mathbb{1} \tag{3.29}
\end{equation*}
$$

Reinstating for clarity the indices $A$ and $B$ the defining relation is

$$
\begin{equation*}
\left(\gamma^{\mu}\right)^{A}{ }_{B}\left(\gamma^{\nu}\right)^{B}{ }_{C}+\left(\gamma^{\nu}\right)^{A}{ }_{B}\left(\gamma^{\mu}\right)^{B}{ }_{C}=2 \eta^{\mu \nu} \mathbb{1}_{C}^{A} . \tag{3.30}
\end{equation*}
$$

This implies

$$
\gamma^{\mu} \gamma^{v}=\left\{\begin{array}{ll}
\eta^{\mu v} & \text { if } \mu=v  \tag{3.31}\\
-\gamma^{v} \gamma^{\mu} & \text { if } \mu \neq v
\end{array} .\right.
$$

Therefore

$$
\begin{equation*}
\left(\gamma^{0}\right)^{2}=\mathbb{1}, \quad\left(\gamma^{i}\right)^{2}=-\mathbb{1} \tag{3.32}
\end{equation*}
$$

The central point is that given $\left(\gamma^{\mu}\right)^{A}{ }_{B}$ as above the objects

$$
\begin{equation*}
\left(S^{\rho \sigma}\right)_{B}^{A}:=\frac{i}{4}\left[\gamma^{\rho}, \gamma^{\sigma}\right]_{B}^{A} \tag{3.33}
\end{equation*}
$$

form a representation of $\operatorname{so}(1,3)$, i.e.

$$
\begin{equation*}
\left[S^{\rho \sigma}, S^{\tau k}\right]=-i\left(\eta^{\rho k} S^{\sigma \tau}+\eta^{\sigma \tau} S^{\rho k}-\eta^{\rho \tau} S^{\sigma k}-\eta^{\sigma k} S^{\rho \tau}\right) \tag{3.34}
\end{equation*}
$$

therefore (3.18) holds. Thus we have constructed a representation of so(1,3) from the Clifford algebra:

One can prove this by direct computation. To this end note that

$$
S^{\mu \nu}=\left\{\begin{array}{cc}
0 & \text { if } \mu=v  \tag{3.35}\\
\frac{i}{2} \gamma^{\mu} \gamma^{\nu} & \text { if } \mu \neq v
\end{array}\right\}=\frac{i}{2}\left(\gamma^{\mu} \gamma^{\nu}-\eta^{\mu \nu}\right) .
$$

The claim then follows with the help of the anti-commutation relations (3.29) after some algebra as worked out in the tutorial.

Now we want to find an explicit representation of the Clifford algebra. Since it is useful to know the result for an arbitrary number of spacetime dimensions let us give the general result valid for Cliff( $1, d-1$ ). The problem is therefore to find $n \times n$-matrices

$$
\begin{equation*}
\left(\gamma^{\mu}\right)^{A}{ }_{B} \text { with } \mu=0,1,2, \ldots, d-1 \text { and } A, B=1, \ldots, n \tag{3.36}
\end{equation*}
$$

subject to (3.29). This will then also give a representation of the Lorentz algebra so $(1, d-1)$. In the tutorials we will prove the following famous theorem:

The irreducible representations of $\operatorname{Cliff}(1, d-1)$ are of dimension

$$
\begin{align*}
& n=\quad 2^{\frac{d}{2}} \text { if dis even } \\
& n=2^{\frac{1}{2}(d-1)} \text { if dis odd } . \tag{3.37}
\end{align*}
$$

Let us now specialise to $d=4$ as is relevant for QFT in four spacetime dimensions. In this case $n=4$. One choice of $\gamma^{\mu}$ called chiral or Dirac representation is

$$
\gamma^{0}=\left(\begin{array}{cc}
0 & \mathbb{1}_{2}  \tag{3.38}\\
\mathbb{1}_{2} & 0
\end{array}\right), \gamma^{i}=\left(\begin{array}{cc}
0 & \sigma^{i} \\
-\sigma^{i} & 0
\end{array}\right),
$$

where $\sigma^{i}$ are the Pauli matrices

$$
\sigma^{1}=\left(\begin{array}{ll}
0 & 1  \tag{3.39}\\
1 & 0
\end{array}\right), \quad \sigma^{2}=\left(\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right), \quad \sigma^{3}=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

with

$$
\begin{equation*}
\left\{\sigma^{i}, \sigma^{j}\right\}=2 \delta^{i j} \tag{3.40}
\end{equation*}
$$

Given an invertible matrix $U$, every other choice $U \gamma^{\mu} U^{-1}$ is also a representation of $\operatorname{Cliff}(1,3)$. The complex vector space on which $\left(\gamma^{\mu}\right)^{A}{ }_{B}$ acts is called the space of Dirac spinors $\psi^{A}, A=1, \ldots, n$ with $n=4$ for $d=4$.
Now, from the above it is clear that the $\psi^{A}$ also form a representation of so $(1,3)$ because the $\gamma^{\mu}$ induce a representation of $s o(1,3)$. More precisely, a Dirac spinor $\psi^{A}$ transforms under a Lorentz transformation

$$
\begin{equation*}
\Lambda_{v}^{\mu}=\left[e^{-\frac{i}{2} \omega_{\rho \sigma} \mathcal{J}^{\rho \sigma}}\right]_{v}^{\mu} \tag{3.41}
\end{equation*}
$$

as

$$
\begin{equation*}
\psi^{A} \mapsto\left[e^{-\frac{i}{2} \omega_{\rho \sigma} S^{\rho \sigma}}\right]_{B}^{A} \psi^{B} . \tag{3.42}
\end{equation*}
$$

A Dirac spinor field $\psi^{A}(x)$ transforms as

$$
\begin{align*}
x & \mapsto x^{\prime}=\Lambda x  \tag{3.43}\\
\psi^{A}(x) & \mapsto[S(\Lambda)]_{B}^{A} \psi^{B}\left(\Lambda^{-1} x^{\prime}\right)=[S(\Lambda)]_{B}^{A} \psi^{B}(x) \tag{3.44}
\end{align*}
$$

with

$$
\begin{equation*}
[S(\Lambda)]_{B}^{A}=\left[e^{-\frac{i}{2} \omega_{\rho \sigma} S^{\rho \sigma}}\right]_{B}^{A} \tag{3.45}
\end{equation*}
$$

The indices $A=1, \ldots, n$, with $n=4$ for $\operatorname{Cliff}(1,3)$ are called spinor indices. Even though in 4 dimensions $A$ happens to run from $1, \ldots, 4$, we must not cunfuse them at any time with spacetime indices $\mu, \nu$.

Now that we found a new type of representation of the Lorentz algebra so(1,3) we would like to give it a physical interpretation. The important claim is that

## A Dirac spinor field $\psi^{A}(x)$ behaves like a spin $\frac{1}{2}$-field.

To prove this we consider a spatial rotation around $\vec{n}$ with

$$
\begin{equation*}
\omega_{i j}=\alpha \epsilon_{i j k} n^{k} \tag{3.46}
\end{equation*}
$$

and

$$
\begin{equation*}
S[\Lambda]=\left[e^{-\frac{i}{2} \omega_{i j} S^{i j}}\right]_{B}^{A} \tag{3.47}
\end{equation*}
$$

Now,

$$
S^{i j}=\frac{i}{4}\left[\gamma^{i}, \gamma^{j}\right]=\left(\begin{array}{cc}
-\frac{i}{2} \sigma^{i} \sigma^{j} & 0  \tag{3.48}\\
0 & -\frac{i}{2} \sigma^{i} \sigma^{j}
\end{array}\right)=\left(\begin{array}{cc}
\frac{1}{2} \sigma^{k} & 0 \\
0 & \frac{1}{2} \sigma^{k}
\end{array}\right)
$$

with ( $i j k$ ) cyclic permutations of 123 . The latter equations follows from the properties of the Pauli matrices. So

$$
S[\Lambda]=\left(\begin{array}{cc}
e^{-i \frac{\alpha}{2} \vec{n} \vec{\sigma}} & 0  \tag{3.49}\\
0 & e^{-i \frac{\alpha}{2} \vec{n} \vec{\sigma}}
\end{array}\right)
$$

What is crucial is the factor of $\frac{1}{2}$, which indeed indicates a transformation as a spin $\frac{1}{2}$ under the subalgebra $\operatorname{so}(3) \subset \operatorname{so}(1,3)$ of spatial rotations. In particular consider a rotation by $\alpha=2 \pi$ around $\vec{n}^{T}=(0,0,1)$,

$$
S[\Lambda]=\exp \left(\begin{array}{cc}
-i \pi \sigma_{3} & 0  \tag{3.50}\\
0 & -i \pi \sigma_{3}
\end{array}\right)=-\mathbb{1}
$$

Therefore

$$
\begin{equation*}
\psi^{A} \mapsto-\psi^{A} \tag{3.51}
\end{equation*}
$$

under a rotation by $2 \pi$.

While we have found a representation of the Lorentz-algebra so( 1,3 ), this does not give a representation of the Lorentz-group $S O(1,3)$, but rather of its double cover $\operatorname{Spin}(1,3)$. The latter is isomorphic to $S L(2, \mathbb{C})$, the group of complex $2 \times 2$ matrices with determinant 1 (not to confuse with $S U(2)$, the group of complex unitary such matrices). In the tutorial we will investigate this isomorphism and the corresponding description of spinors.

A Dirac spinor forms a representation of $\operatorname{Spin}(1,3) \simeq S L(2, \mathbb{C})$, not of $S O(1,3)$.

The reason is that $S[\Lambda=\mathbb{1}] \neq \mathbb{1}$, which is incompatible with the group law of $S O(1,3)$. In $\operatorname{Spin}(1,3)$ by contrast, a rotation around $2 \pi$ does not correspond to the $\mathbb{1}$ because it is defined as the double cover of $S O(1,3) .^{3}$ This is the relativistic version of the statement familiar from Quantum Mechanics that the $j=\frac{1}{2}$ spinor representation of the algebra of spatial rotations so(3) does not furnish a representation of the Lie group $S O(3)$ but only of its double cover $S U(2)$.

### 3.3 The Dirac action

We now want to construct Lorentz scalars and vectors out of $\psi^{A}(x)$ in order to construct a covariant action. To do so we define the conjugate spinor $\psi^{\dagger}:=\left(\psi^{*}\right)^{T}$, i.e.

$$
\begin{equation*}
\psi^{\dagger}=\left(\left(\psi^{1}\right)^{*},\left(\psi^{2}\right)^{*},\left(\psi^{3}\right)^{*},\left(\psi^{4}\right)^{*}\right) . \tag{3.52}
\end{equation*}
$$

We now would like to know how $\psi^{\dagger}(x) \psi(x)$ behaves under a Lorentz transformation $\Lambda$, under which

$$
\begin{align*}
x & \mapsto x^{\prime}=\Lambda x  \tag{3.53}\\
\psi(x) & \mapsto S(\Lambda) \psi\left(\Lambda^{-1} x^{\prime}\right)  \tag{3.54}\\
\psi^{\dagger}(x) & \mapsto \psi^{\dagger}\left(\Lambda^{-1} x^{\prime}\right) S^{\dagger}(\Lambda) \tag{3.55}
\end{align*}
$$

In order for $\psi^{\dagger}(x) \psi(x)$ to transform as a Lorentz scalar, we would need $S(\Lambda)^{\dagger}=S^{-1}(\Lambda)$. However,

$$
\begin{equation*}
S^{\dagger}(\Lambda)=e^{\left[-\frac{i}{2} \omega_{\rho \sigma} S^{\rho \sigma}\right]^{\dagger}}=e^{\frac{i}{2} \omega_{\rho \sigma}\left(S^{\rho \sigma}\right)^{\dagger}} \tag{3.56}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(S^{\rho \sigma}\right)^{\dagger}=\left(\frac{i}{4}\left[\gamma^{\rho}, \gamma^{\sigma}\right]\right)^{\dagger}=-\frac{i}{4}\left[\gamma^{\sigma \dagger}, \gamma^{\rho \dagger}\right]=\frac{i}{4}\left[\gamma^{\rho^{\dagger}}, \gamma^{\sigma \dagger}\right] . \tag{3.57}
\end{equation*}
$$

It is not possible to pick all $\gamma^{\mu}$ hermitian at the same time since

$$
\begin{gather*}
\left(\gamma^{0}\right)^{2}=1 \rightarrow \text { real eigenvalues, } \\
\left(\gamma^{i}\right)^{2}=-1 \rightarrow \text { imaginary eigenvalues. } \tag{3.58}
\end{gather*}
$$

Therefore, $\left(S^{\rho \sigma}\right)^{\dagger} \neq\left(S^{\rho \sigma}\right)$ and thus $S(\Lambda)^{\dagger} \neq S^{-1}(\Lambda)$.

[^17]Rather, it is possible to pick $\gamma^{\mu}$ such that

$$
\begin{equation*}
\left(\gamma^{0}\right)^{\dagger}=\gamma^{0}, \quad\left(\gamma^{i}\right)^{\dagger}=-\gamma^{i}, \tag{3.59}
\end{equation*}
$$

which is indeed satisfied by the Dirac representation above. In this case

$$
\begin{equation*}
\left(\gamma^{\mu}\right)^{\dagger}=\gamma^{0} \gamma^{\mu} \gamma^{0}, \tag{3.60}
\end{equation*}
$$

because $\gamma^{0} \gamma^{0}=\mathbb{1}$ and $\gamma^{0} \gamma^{i}=-\gamma^{i} \gamma^{0}$. Therefore

$$
\begin{equation*}
\gamma^{0}\left(S^{\rho \sigma}\right)^{\dagger} \gamma^{0}=\frac{i}{4} \gamma^{0}\left[\gamma^{\rho \dagger}, \gamma^{\sigma \dagger}\right] \gamma^{0}=\frac{i}{4}\left[\gamma^{\rho}, \gamma^{\sigma}\right]=S^{\rho \sigma} \tag{3.61}
\end{equation*}
$$

and so

$$
\begin{equation*}
\gamma^{0} S[\Lambda]^{\dagger} \gamma^{0}=S[\Lambda]^{-1} \tag{3.62}
\end{equation*}
$$

Inspired by this result we define the Dirac conjugate spinor

$$
\begin{equation*}
\bar{\psi}:=\psi^{\dagger} \gamma^{0} . \tag{3.63}
\end{equation*}
$$

It transforms as

$$
\begin{equation*}
\bar{\psi}(x) \mapsto \psi^{\dagger}\left(\Lambda^{-1} x^{\prime}\right) S^{\dagger}(\Lambda) \gamma^{0}=\psi^{\dagger}\left(\Lambda^{-1} x^{\prime}\right) \gamma^{0} S^{-1}(\Lambda)=\bar{\psi}\left(\Lambda^{-1} x^{\prime}\right) S^{-1}(\Lambda) \tag{3.64}
\end{equation*}
$$

Thus the spinor bilinear $\bar{\psi}(x) \psi(x)$ transforms as

$$
\begin{equation*}
\bar{\psi}(x) \psi(x) \mapsto \bar{\psi}\left(\Lambda^{-1} x^{\prime}\right) S^{-1}(\Lambda) S(\Lambda) \psi\left(\left(\Lambda^{-1} x^{\prime}\right)\right)=\bar{\psi}\left(\Lambda^{-1} x^{\prime}\right) \psi\left(\Lambda^{-1} x^{\prime}\right) \tag{3.65}
\end{equation*}
$$

So $\bar{\psi}(x) \psi(x)$ is a scalar quantity. Furthermore one can show that $\bar{\psi}(x) \gamma^{\mu} \psi(x)$ transforms as a vector because

$$
\begin{equation*}
\bar{\psi}(x) \gamma^{\mu} \psi(x) \mapsto \bar{\psi}\left(\Lambda^{-1} x^{\prime}\right) S^{-1}(\Lambda) \gamma^{\mu} S(\Lambda) \psi\left(\Lambda^{-1} x^{\prime}\right) \tag{3.66}
\end{equation*}
$$

and

$$
\begin{equation*}
S^{-1}(\Lambda) \gamma^{\mu} S(\Lambda)=\Lambda_{v}^{\mu} \gamma^{v} \tag{3.67}
\end{equation*}
$$

as will be proven in the tutorials. This logic can be repeated for tensors of rank $n$.

We can now build an action for $\psi(x)$, where we follow the principle of simplicity. The simplest Lorentz scalars built from $\psi(x)$ which include non-trivial dynamics (i.e. at least one derivative has to appear) are

$$
\begin{equation*}
\bar{\psi}(x) \gamma^{\mu} \partial_{\mu} \psi(x), \quad \bar{\psi}(x) \psi(x) \tag{3.68}
\end{equation*}
$$

So we take as the action for the free classical Dirac spinor field

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x \bar{\psi}(x)\left[i \gamma^{\mu} \partial_{\mu}-m\right] \psi(x) \tag{3.69}
\end{equation*}
$$

where

- the factor $i$ is required for the action to be real
- and $|m|$ will be the mass of the Dirac spinor particle.

This is the simplest action we can build. Interestingly - unlike in for the scalar field - it is possible to construct a Lorentz invariant action with only one derivative. Furthermore note that $\psi(x)$ has mass dimension [mass] ${ }^{3 / 2}$.

Since $\psi(x)$ is complex we treat $\psi(x)$ and $\psi^{\dagger}(x)$ as independent when deriving the equations of motion. Varying with respect to $\psi^{\dagger}(x)$ yields the Dirac equation

$$
\begin{equation*}
\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi(x)=0, \tag{3.70}
\end{equation*}
$$

which is linear in the derivatives. In order to see that this indeed describes a field with mass $|m|$ we note that

$$
\begin{align*}
0 & =\left(i \gamma^{\mu} \partial_{\mu}+m\right)\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi \\
& =\left(-\gamma^{\mu} \gamma^{\nu} \partial_{\mu} \partial_{\nu}-m^{2}\right) \psi \\
& =-(\underbrace{\frac{1}{2}\left\{\gamma^{\mu}, \gamma^{\nu}\right\}}_{=\eta^{\mu \nu}} \partial_{\mu} \partial_{\nu}+m^{2}) \psi \tag{3.71}
\end{align*}
$$

and thus obtain the Klein-Gordon equation (1.29) for the spinor field $\psi(x)$. Loosely speaking:

$$
\begin{equation*}
\text { Dirac equation }=\sqrt{\text { Klein-Gordon equation. }} \tag{3.72}
\end{equation*}
$$

This is a consequence of the particular manner how we constructed a representation of the Lorentz algebra from the Clifford algebra. Equ. (3.71) also justifies the relative factor of $i$ in the Dirac action, for which we can take $m$ to be real.

### 3.4 Chirality and Weyl spinors

The Dirac spinor representation of $\operatorname{Cliff}(1,3)$ is not irreducible as a representation of $\operatorname{Spin}(1,3)$. Indeed for our special choice of Dirac matrices

$$
\gamma^{0}=\left(\begin{array}{cc}
0 & \mathbb{1}_{2}  \tag{3.73}\\
\mathbb{1}_{2} & 0
\end{array}\right), \quad \gamma^{i}=\left(\begin{array}{cc}
0 & \sigma^{i} \\
-\sigma^{i} & 0
\end{array}\right)
$$

$\psi^{A}$ transforms under a spatial rotation as

$$
\psi \mapsto\left(\begin{array}{cc}
e^{-\frac{i}{2} \alpha \vec{n} \cdot \vec{\sigma}} & 0  \tag{3.74}\\
0 & e^{-\frac{i}{2} \alpha \vec{k} \cdot \vec{\sigma}}
\end{array}\right) \psi
$$

and under a Lorentz boost with $\omega_{0 j}=\chi_{j}$ as

$$
\psi \mapsto\left(\begin{array}{cc}
e^{-\frac{1}{2} \vec{\gamma} \cdot \vec{\sigma}} & 0  \tag{3.75}\\
0 & e^{+\frac{1}{2} \vec{\gamma} \cdot \vec{\sigma}}
\end{array}\right) \psi .
$$

Therefore the subspaces spanned by

$$
\begin{equation*}
\psi_{-}^{T}=\left(\psi^{1}, \psi^{2}, 0,0\right) \text { and } \psi_{+}^{T}=\left(0,0, \psi^{3}, \psi^{4}\right) \tag{3.76}
\end{equation*}
$$

transform separately. Irrespective of the concrete representation the reducibility of the Dirac spinor repsentation as a representation of $\operatorname{Spin}(1,3)$ can be seen as follows: Define

$$
\begin{equation*}
\gamma^{5}=i \gamma^{0} \gamma^{1} \gamma^{2} \gamma^{3} \tag{3.77}
\end{equation*}
$$

with the properties

$$
\begin{equation*}
\left(\gamma^{5}\right)^{2}=\mathbb{1}, \quad\left\{\gamma^{5}, \gamma^{\mu}\right\}=0, \quad\left[S^{\rho \sigma}, \gamma^{5}\right]=0 \tag{3.78}
\end{equation*}
$$

following from the properties of the Clifford-algebra. Now consider the orthogonal projection operators

$$
\begin{equation*}
\mathbb{P}_{ \pm}:=\frac{1}{2}\left(\mathbb{1} \pm \gamma^{5}\right) \tag{3.79}
\end{equation*}
$$

with the properties

$$
\begin{equation*}
\left(\mathbb{P}_{ \pm}\right)^{2}=\mathbb{P}_{ \pm}, \quad \mathbb{P}_{+} \mathbb{P}_{-}=0=\mathbb{P}_{-} \mathbb{P}_{+}, \quad \mathbb{1}=\mathbb{P}_{+}+\mathbb{P}_{-} \tag{3.80}
\end{equation*}
$$

Defining

$$
\begin{equation*}
\psi_{ \pm}:=\mathbb{P}_{ \pm} \psi \tag{3.81}
\end{equation*}
$$

yields

$$
\begin{equation*}
\mathbb{P}_{ \pm} \psi_{\mp}=0 . \tag{3.82}
\end{equation*}
$$

Now, since $\left[S^{\rho \sigma}, \gamma^{5}\right]=0$ we have $\mathbb{P}_{\mp} S[\Lambda] \psi_{ \pm}=0$, i.e. the $\pm$ subspaces transform separately under $\operatorname{Spin}(1,3)$. The $\psi_{ \pm}$are called positive/negative-chirality spinors.
In the special representation

$$
\gamma^{0}=\left(\begin{array}{cc}
0 & \mathbb{1}_{2}  \tag{3.83}\\
\mathbb{1}_{2} & 0
\end{array}\right), \quad \gamma^{i}=\left(\begin{array}{cc}
0 & \sigma^{i} \\
-\sigma^{i} & 0
\end{array}\right), \quad \gamma^{5}=\left(\begin{array}{cc}
-\mathbb{1}_{2} & 0 \\
0 & \mathbb{1}_{2}
\end{array}\right)
$$

we find

$$
\begin{equation*}
\psi_{-}=\binom{u_{-}}{0}, \quad \psi_{+}=\binom{0}{u_{+}} . \tag{3.84}
\end{equation*}
$$

The 2-component objects $u_{ \pm}$are called Weyl-spinors. It is the Weyl spinors that form irreducible representations of $\operatorname{Spin}(1,3)$.

It is instructive to rewrite the Dirac action in this language. The Dirac action decomposes as follows: Since $\gamma^{0} \gamma^{5}=-\gamma^{5} \gamma^{0}$ we have

$$
\begin{equation*}
\gamma^{0} \mathbb{P}_{ \pm}=\mathbb{P}_{\mp} \gamma^{0} \tag{3.85}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\overline{\psi_{ \pm}}=\left(\mathbb{P}_{ \pm} \psi\right)^{\dagger} \gamma^{0}=\psi^{\dagger} \mathbb{P}_{ \pm} \gamma^{0}=\bar{\psi} \mathbb{P}_{\mp} \tag{3.86}
\end{equation*}
$$

This means

$$
\begin{align*}
& \overline{\psi_{ \pm}} \psi_{ \pm}=\bar{\psi} \mathbb{P}_{\mp} \mathbb{P}_{ \pm} \psi=0,  \tag{3.87}\\
& \overline{\psi_{ \pm}} \psi_{\mp}=\bar{\psi} \mathbb{P}_{\mp} \psi \neq 0
\end{align*}
$$

and similary (using in addition that $\gamma^{\mu} \mathbb{P}_{ \pm}=\mathbb{P}_{\mp} \gamma^{\mu}$ )

$$
\begin{align*}
& \overline{\psi_{ \pm}} \gamma^{\mu} \psi_{ \pm}=\bar{\psi} \gamma^{\mu} \mathbb{P}_{ \pm} \psi \neq 0 \\
& \overline{\psi_{\mp}} \gamma^{\mu} \psi_{ \pm}=0 . \tag{3.88}
\end{align*}
$$

The Dirac action can then be written as

$$
\begin{align*}
S & =\int \mathrm{d}^{4} x \bar{\psi}\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi  \tag{3.89}\\
& =\int \mathrm{d}^{4} x\left[\overline{\psi_{+}} i \gamma^{\mu} \partial_{\mu} \psi_{+}+\overline{\psi_{-}} i \gamma^{\mu} \partial_{\mu} \psi_{-}-m\left(\overline{\psi_{+}} \psi_{-}+\overline{\psi_{-}} \psi_{+}\right)\right],
\end{align*}
$$

or in Weyl-spinor notation as

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left[u_{+}^{\dagger} i \sigma^{\mu} \partial_{\mu} u_{+}+u_{-}^{\dagger} i \bar{\sigma}^{\mu} \partial_{\mu} u_{-}-m\left(u_{+}^{\dagger} u_{-}+u_{-}^{\dagger} u_{+}\right)\right] . \tag{3.90}
\end{equation*}
$$

Here

$$
\begin{equation*}
\sigma^{\mu} \equiv\left(\mathbb{1}_{2}, \sigma^{i}\right) \text { and } \bar{\sigma}^{\mu}=\left(\mathbb{1}_{2},-\sigma^{i}\right) . \tag{3.91}
\end{equation*}
$$

From the decomposed Dirac equation we can draw the following important conclusions about the underlying physics:

- If $m=0, u_{+}$and $u_{-}$decouple and describe independent degrees of freedom subject to the Weyl equations

$$
\begin{equation*}
i \sigma^{\mu} \partial_{\mu} u_{+}(x)=0, \quad i \bar{\sigma}^{\mu} \partial_{\mu} u_{-}(x)=0 \tag{3.92}
\end{equation*}
$$

Both $u_{+}$and $u_{-}$transform in the $s=\frac{1}{2}$ representation of $\operatorname{SU}(2) \subset \operatorname{Spin}(1,3)$, in the sense that under spatial rotations

$$
\begin{align*}
x & \mapsto R x,  \tag{3.93}\\
u_{ \pm}(x) & \mapsto e^{-\frac{i}{2} \alpha \vec{n} \cdot \vec{\sigma}} u_{ \pm}\left(R^{-1} x\right) . \tag{3.94}
\end{align*}
$$

We define the helicity

$$
\begin{equation*}
h=\epsilon_{i j k} \hat{p}^{i} S^{j k}=\frac{1}{2} \hat{\vec{p}} \cdot \vec{\sigma}, \tag{3.95}
\end{equation*}
$$

which is the projection of the spin onto the momentum direction $\hat{\vec{p}}$. To solve the Weyl equations of momentum $p^{\mu}$ one makes the ansatz (see the next chapter for details)

$$
\begin{equation*}
u_{ \pm}(x)=u_{ \pm}(p) e^{-i p \cdot x} \tag{3.96}
\end{equation*}
$$

and finds

$$
\begin{equation*}
h u_{ \pm}(p)= \pm \frac{1}{2} u_{ \pm}(p), \tag{3.97}
\end{equation*}
$$

where $u_{+}$corresponds to right-handed and $u_{-}$to left-handed spinors.

- If $m \neq 0, u_{+}$and $u_{-}$do not decouple. In this sense the full 4-component Dirac spinor $\psi$ is needed to describe massive spin $-\frac{1}{2}$ fields. Note for fields of mass $m \neq 0$ it is impossible to define a Lorentz invariant notion of helicity because massive particles travel with a velocity $v<c$ and it is always possible to find a Lorentz frame in which the particle moves into the opposite direction. This causes a change in the helicity.


### 3.5 Classical plane-wave solutions

To solve the classical equation of motion (3.70) we make the ansatz

$$
\begin{equation*}
\psi(x)=u(\vec{p}) e^{-i p \cdot x} \tag{3.98}
\end{equation*}
$$

with $p=\left(E_{p}, \vec{p}\right)$ and $E_{p}=\sqrt{\vec{p}^{2}+m^{2}}$. Here we have used that the dispersion relation $p^{2}-m^{2}=0$ is satisfied because $\psi(x)$ obeys the Klein-Gordon equation (3.71). Plugging this ansatz into the Dirac equation yields

$$
\begin{equation*}
\left(p^{\mu} \gamma_{\mu}-m\right) u(\vec{p})=0 . \tag{3.99}
\end{equation*}
$$

For our choice of $\gamma^{\mu}$ this is

$$
\left(\begin{array}{cc}
-m & p_{\mu} \sigma^{\mu}  \tag{3.100}\\
p_{\mu} \bar{\sigma}^{\mu} & -m
\end{array}\right) u(\vec{p})=0
$$

We make the ansatz

$$
\begin{equation*}
u(\vec{p})=\binom{u_{1}}{u_{2}}, \tag{3.101}
\end{equation*}
$$

which leads to the relations

$$
\begin{align*}
p_{\mu} \sigma^{\mu} u_{2} & =m u_{1},  \tag{3.102}\\
p_{\mu} \bar{\sigma}^{\mu} u_{1} & =m u_{2} .
\end{align*}
$$

These are consistent because

$$
\begin{equation*}
\left(p_{\mu} \sigma^{\mu}\right)\left(p_{\mu} \bar{\sigma}^{\mu}\right)=\left(p^{0}\right)^{2}-p_{i} p_{j} \underbrace{\sigma^{i} \sigma^{j}}_{\delta_{i j}}=m^{2} . \tag{3.103}
\end{equation*}
$$

The general solution can be parametrised by introducing some 2 -component Weyl spinor $\xi^{\prime}$ and writing

$$
\begin{equation*}
u(\vec{p})=\binom{\frac{1}{m} p \cdot \sigma \xi^{\prime}}{\xi^{\prime}} \tag{3.104}
\end{equation*}
$$

We make a conventional choice

$$
\begin{equation*}
\xi^{\prime}=\sqrt{p \cdot \bar{\sigma}} \xi \tag{3.105}
\end{equation*}
$$

with $\xi^{\dagger} \xi=1$. Here $\sqrt{M}$ denotes the matrix whose eigenvalues are the square root of those of $M$. With this choice we find the so-called positive frequency solution

$$
\begin{equation*}
\psi(x)=u(\vec{p}) e^{-i p \cdot x}, \quad u(\vec{p})=\binom{\sqrt{p \cdot \sigma} \xi}{\sqrt{p \cdot \bar{\sigma}} \xi} \tag{3.106}
\end{equation*}
$$

Likewise

$$
\begin{equation*}
\psi(x)=v(\vec{p}) e^{i p \cdot x} \tag{3.107}
\end{equation*}
$$

is a solution if

$$
\begin{equation*}
(\gamma \cdot p+m) v(\vec{p})=0 . \tag{3.108}
\end{equation*}
$$

This yields the negative frequency solution

$$
\begin{equation*}
\psi(x)=v(\vec{p}) e^{i p \cdot x}, \quad v(\vec{p})=\binom{\sqrt{p \cdot \sigma} \xi}{-\sqrt{p \cdot \bar{\sigma}} \xi} \tag{3.109}
\end{equation*}
$$

Let us introduce a basis of the space of 2 -spinors by $\xi_{s}$ with

$$
\begin{equation*}
\xi_{\frac{1}{2}}=\binom{1}{0}, \quad \xi_{-\frac{1}{2}}=\binom{0}{1} \tag{3.110}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\xi_{s}^{\dagger} \xi_{s^{\prime}}=\delta_{s s^{\prime}} \tag{3.111}
\end{equation*}
$$

Note that $\xi_{s}$, viewed as a Weyl spinor, transforms in the $s=\frac{1}{2}$ representation of $\operatorname{SU}(2)$. In particular

$$
\begin{equation*}
\frac{1}{2} \sigma_{3} \xi_{ \pm \frac{1}{2}}= \pm \frac{1}{2} \xi_{ \pm \frac{1}{2}} . \tag{3.112}
\end{equation*}
$$

Correspondingly $u_{s}$ and $v_{s}$ with $s= \pm \frac{1}{2}$ describe spinors with spin $\pm \frac{1}{2}$ in direction $x_{3}$. In the tutorial we will convince ourselves of the important identities

$$
\begin{align*}
& \bar{u}_{s}(\vec{p}) u_{s^{\prime}}(\vec{p})=2 m \delta_{s s^{\prime}}, \quad u_{s}^{\dagger}(\vec{p}) u_{s^{\prime}}(\vec{p})=2 p^{0} \delta_{s s^{\prime}} \\
& \bar{v}_{s}(\vec{p}) v_{s^{\prime}}(\vec{p})=-2 m \delta_{s s^{\prime}}, \quad v_{s}^{\dagger}(\vec{p}) v_{s^{\prime}}(\vec{p})=2 p^{0} \delta_{s s^{\prime}},  \tag{3.113}\\
& \bar{u}_{s}(\vec{p}) v_{s^{\prime}}(\vec{p})=0
\end{align*}
$$

and

$$
\begin{align*}
& \sum_{s} u_{s}(\vec{p}) \bar{u}_{s}(\vec{p})=\gamma \cdot p+m  \tag{3.114}\\
& \sum_{s} v_{s}(\vec{p}) \bar{v}_{s}(\vec{p})=\gamma \cdot p-m
\end{align*}
$$

Let us consider the following example: Suppose we have chosen coordinates such that $p^{\mu}=\left(E, 0,0, p^{3}\right)$. According to the above, the Dirac spinor solution $u_{\frac{1}{2}}(\vec{p})$ with spin $\frac{1}{2}$ along $x_{3}$ is given by

$$
\begin{equation*}
u_{\frac{1}{2}}(\vec{p})=\binom{\sqrt{\sigma \cdot p} \xi_{\frac{1}{2}}}{\sqrt{\bar{\sigma} \cdot p} \xi_{\frac{1}{2}}}=\binom{\sqrt{E-p^{3}} \xi_{\frac{1}{2}}}{\sqrt{E+p^{3}} \xi_{\frac{1}{2}}} . \tag{3.115}
\end{equation*}
$$

If $m=0$, then $E=p^{3}$ and

$$
u_{\frac{1}{2}}(\vec{p})=\sqrt{2 E}\left(\begin{array}{l}
0  \tag{3.116}\\
0 \\
1 \\
0
\end{array}\right) \equiv u_{+}(\vec{p})
$$

since $\xi_{\frac{1}{2}}=\binom{1}{0}$. The solution correspnding to $\xi_{-\frac{1}{2}}=\binom{0}{1}$ is

$$
u_{-\frac{1}{2}}(\vec{p})=\binom{\sqrt{E+p^{3}} \xi_{-\frac{1}{2}}}{\sqrt{E-p^{3}} \xi_{-\frac{1}{2}}} \xrightarrow{m=0} \sqrt{2 E}\left(\begin{array}{l}
0  \tag{3.117}\\
1 \\
0 \\
0
\end{array}\right) \equiv u_{-}(\vec{p})
$$

The helicity of the solutions is

$$
\begin{align*}
h u_{+}(\vec{p}) & =\frac{1}{2} u_{+}(\vec{p}) \Rightarrow \text { right-handed } \\
h u_{-}(\vec{p}) & =-\frac{1}{2} u_{-}(\vec{p}) \Rightarrow \text { left-handed. } \tag{3.118}
\end{align*}
$$

### 3.6 Quantisation of the Dirac field

So far our analysis of spinors has been classical. To define the quantum theory of spinor fields we follow the same procedure as in the scalar case, but we will encounter a problem if we just impose the same commutation relation as in the scalar case. Its resolution will prove the fermionic nature of spin $\frac{1}{2}$ particles.

### 3.6.1 Using the commutator

Starting from the classical Lagrangian

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x \bar{\psi}\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi \tag{3.119}
\end{equation*}
$$

we find the conjugate momentum density

$$
\begin{equation*}
\Pi_{A}=\frac{\partial \mathcal{L}}{\partial \dot{\psi}^{A}}=i \psi_{A}^{\dagger} \tag{3.120}
\end{equation*}
$$

We promote $\psi(x)$ and $\Pi(x)$ to Schrödinger picture quantum operators and expand these as

$$
\begin{align*}
\psi(\vec{x}) & =\sum_{s} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left[a_{s}(\vec{p}) u_{s}(\vec{p}) e^{i \vec{p} \cdot \vec{x}}+b_{s}^{\dagger}(\vec{p}) v_{s}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right],  \tag{3.121}\\
\psi^{\dagger}(\vec{x}) & =\sum_{s} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left[b_{s}(\vec{p}) v_{s}^{\dagger}(\vec{p}) e^{i \vec{p} \cdot \vec{x}}+a_{s}^{\dagger}(\vec{p}) u_{s}^{\dagger}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right],
\end{align*}
$$

where $u_{s}$ and $v_{s}$ are the spinor-valued solutions to the classical equation (3.70) corresponding to the ansatz (3.98). The quantum operators $a_{s}$ and $b_{s}$ are independent because the field $\psi$ is complex. This expansion guarantees that the corresponding Heisenberg fields with $e^{ \pm i \vec{p} \cdot \vec{x}} \mapsto e^{\mp i p \cdot x}$ satisfy the Dirac equation as an operator equation.

We are now careful with operator orderings and proceed without specifying the commutation relations. The Hamiltonian density is

$$
\begin{equation*}
\mathcal{H}=\Pi \dot{\psi}-\mathcal{L}=\bar{\psi}\left(-i \gamma^{j} \partial_{j}+m\right) \psi \tag{3.122}
\end{equation*}
$$

Now,

$$
\begin{align*}
\left(-i \gamma^{j} \partial_{j}+m\right) \psi= & \sum_{s} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left[\left(-\gamma^{j} p_{j}+m\right) u_{s}(\vec{p}) a_{s}(\vec{p}) e^{i \vec{p} \cdot \vec{x}}\right.  \tag{3.123}\\
& \left.+\left(\gamma^{j} p_{j}+m\right) v_{s}(\vec{p}) b_{s}^{\dagger}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right]
\end{align*}
$$

because

$$
\begin{equation*}
\partial_{j} e^{i \vec{p} \cdot \vec{x}}=\partial_{j} e^{-i p_{k} x^{k}}=-i p_{j} e^{i \vec{p} \cdot \vec{x}} \tag{3.124}
\end{equation*}
$$

Now use that $v_{s}$ and $u_{s}$ solve the classical equation of motion, i.e.

$$
\begin{align*}
& \left(\gamma^{\mu} p_{\mu}-m\right) u_{s}(\vec{p})=0 \Rightarrow\left(-\gamma^{j} p_{j}+m\right) u_{s}(\vec{p})=\gamma^{0} p_{0} u_{s}(\vec{p}),  \tag{3.125}\\
& \left(\gamma^{\mu} p_{\mu}+m\right) v_{s}(\vec{p})=0 \Rightarrow\left(\gamma^{j} p_{j}+m\right) v_{s}(\vec{p})=-\gamma^{0} p_{0} v_{s}(\vec{p})
\end{align*}
$$

and therefore

$$
\begin{equation*}
\left(-i \gamma^{j} \partial_{j}+m\right) \psi=\sum_{s} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \sqrt{\frac{E_{p}}{2}} \gamma^{0}\left[a_{s}(\vec{p}) u_{s}(\vec{p}) e^{i \vec{p} \cdot \vec{x}}-b_{s}^{\dagger}(\vec{p}) v_{s}(\vec{p}) e^{-i \vec{p} \cdot \vec{x}}\right] \tag{3.126}
\end{equation*}
$$

To compute the Hamiltonian

$$
\begin{equation*}
H=\int \mathrm{d}^{3} x \psi^{\dagger} \gamma^{0}\left(-i \gamma^{j} \partial_{j}+m\right) \psi \tag{3.127}
\end{equation*}
$$

we furthermore exploit the spinor identities

$$
\begin{align*}
& v_{r}^{\dagger}(\vec{p}) u_{s}(\vec{p})=0, \quad u_{r}^{\dagger}(\vec{p}) v_{s}(\vec{p})=0 \\
& u_{r}^{\dagger}(p) u_{s}(p)=2 p^{0} \delta_{r s}=v_{r}^{\dagger}(p) v_{s}(p) \tag{3.128}
\end{align*}
$$

which eventually give

$$
\begin{equation*}
H=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} E_{p} \sum_{s}\left[a_{s}^{\dagger}(\vec{p}) a_{s}(\vec{p})-b_{s}(\vec{p}) b_{s}^{\dagger}(\vec{p})\right] . \tag{3.129}
\end{equation*}
$$

So far no re-ordering of operators has been performed. Now, the naive guess would be to impose, at the next step, canonical commutation relations as for the scalar field,

$$
\begin{align*}
{\left[\psi^{A}(\vec{x}), \Pi_{B}\left(\vec{x}^{\prime}\right)\right] } & \stackrel{?}{=} i \delta_{B}^{A} \delta^{(3)}\left(\vec{x}-\vec{x}^{\prime}\right) \\
{\left[\psi^{A}(\vec{x}), \psi_{B}^{\dagger}\left(\vec{x}^{\prime}\right)\right] } & \stackrel{?}{=} \delta_{B}^{A} \delta^{(3)}\left(\vec{x}-\vec{x}^{\prime}\right),  \tag{3.130}\\
{\left[\psi^{A}(\vec{x}), \psi^{B}\left(\vec{x}^{\prime}\right)\right] } & \stackrel{?}{=} 0 .
\end{align*}
$$

We will now show why this is wrong. The commutation relations would imply for the modes

$$
\begin{align*}
{\left[a_{r}(\vec{p}), a_{s}^{\dagger}(\vec{q})\right] } & =(2 \pi)^{3}, \delta_{r s}, \delta^{(3)}(\vec{p}-\vec{q}), \\
{\left[b_{r}(\vec{p}), b_{s}^{\dagger}(\vec{q})\right] } & =-(2 \pi)^{3}, \delta_{r s}, \delta^{(3)}(\vec{p}-\vec{q}),  \tag{3.131}\\
{\left[a_{r}(\vec{p}), b_{s}^{(\dagger)}(\vec{q})\right] } & =0 .
\end{align*}
$$

Indeed we have

$$
\begin{align*}
{\left[\psi(\vec{x}), \psi^{\dagger}(\vec{y})\right]=} & \sum_{r, s} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \frac{\mathrm{~d}^{3} q}{(2 \pi)^{3}} \frac{1}{\sqrt{4 E_{p} E_{q}}}\left(\left[a_{r}(\vec{p}), a_{s}^{\dagger}(\vec{q})\right] u_{r}(\vec{p}) u_{s}(\vec{q})^{\dagger} e^{i(\vec{p} \cdot \vec{x}-\vec{q} \cdot \vec{y})}\right.  \tag{3.132}\\
& \left.+\left[b_{r}^{\dagger}(\vec{p}), b_{s}(\vec{q})\right] v_{r}(\vec{p}) v_{s}^{\dagger}(\vec{q}) e^{-i(\vec{p} \cdot \vec{x}-\vec{q} \cdot \vec{y})}\right) .
\end{align*}
$$

Then, using $\sum_{s} u_{s}(\vec{p}) \bar{u}_{s}(\vec{p})=\gamma \cdot p+m$ and $\sum_{s} v_{s}(\vec{p}) \bar{v}_{s}(\vec{p})=\gamma \cdot p-m$, the commutator relations (3.130) of the field $\psi$ would follow from the commutator relations (3.131) of the modes. Note that

$$
\begin{equation*}
\left[a_{r}(\vec{p}), a_{s}^{\dagger}(\vec{q})\right] \text { and }\left[b_{r}^{\dagger}(\vec{p}), b_{s}(\vec{q})\right] \tag{3.133}
\end{equation*}
$$

appear on same footing in the commutator of the field. This ordering is the reason for the crucial relative minus sign in (3.131). If this were correct, then we could reorder the Hamiltonian to find

$$
\begin{equation*}
H=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} E_{p} \sum_{s}\left[a_{s}^{\dagger}(\vec{p}) a_{s}(\vec{p})-b_{s}^{\dagger}(\vec{p}) b_{s}(\vec{p})+(2 \pi)^{3} \delta(0)\right] . \tag{3.134}
\end{equation*}
$$

To give the mode operators a physical interpretation we observe that

$$
\begin{align*}
{\left[H, a_{s}(\vec{p})\right] } & =-E_{p} a_{s}(\vec{p}),  \tag{3.135}\\
{\left[H, a_{s}^{\dagger}(\vec{p})\right] } & =E_{p} a_{s}^{\dagger}(\vec{p})
\end{align*}
$$

and also

$$
\begin{align*}
{\left[H, b_{s}(\vec{p})\right] } & =-E_{p} b_{s}(\vec{p}),  \tag{3.136}\\
{\left[H, b_{s}^{\dagger}(\vec{p})\right] } & =E_{p} b_{s}^{\dagger}(\vec{p})
\end{align*}
$$

because the minus sign in the Hamiltonian cancels with the sign we pick up in the commutator. So as in the scalar theory we would interpret

$$
\begin{align*}
a_{s}(\vec{p}), b_{s}(\vec{p}) & \leftrightarrow \text { annihilation operators }  \tag{3.137}\\
a_{s}^{\dagger}(\vec{p}), b_{s}^{\dagger}(\vec{p}) & \leftrightarrow \text { creation operators } \tag{3.138}
\end{align*}
$$

and define $|0\rangle$ such that $a_{s}(\vec{p})|0\rangle=b_{s}(\vec{p})|0\rangle$. The creation operators give positive energy modes by construction.
While at first sight everything looks fine, in actuality this whole construction is in deep conflict with unitarity. The reason is that due to the minus sign in the commutator

$$
\begin{equation*}
\left[b_{r}(\vec{p}), b_{s}^{\dagger}(\vec{q})\right]=-(2 \pi)^{3} \delta_{r s} \delta^{(3)}(\vec{p}-\vec{q}) \tag{3.139}
\end{equation*}
$$

the $b$-mode excitations are negative norm states. To see this note that

$$
\begin{equation*}
\langle 0|\left[b_{r}(\vec{p}), b_{s}^{\dagger}(\vec{q})\right]|0\rangle=-(2 \pi)^{3} \delta_{r s} \delta^{(3)}(\vec{p}-\vec{q})\langle 0 \mid 0\rangle . \tag{3.140}
\end{equation*}
$$

Evaluated for $r=s$ and $\vec{p}=\vec{q}$ we conclude

$$
\begin{equation*}
0>-(2 \pi)^{3} \delta^{(3)}(0)=\langle 0|\left[b_{r}(\vec{p}), b_{r}^{\dagger}(\vec{p})\right]|0\rangle=\langle 0| b_{r}(\vec{p}) b_{r}^{\dagger}(\vec{p})|0\rangle=\| b_{r}^{\dagger}|0\rangle \|^{2} . \tag{3.141}
\end{equation*}
$$

Therefore no positive norm Hilbert space interpretation is possible!
One might hope to avoid this problem by switching the interpretation of creation versus annihilation operators for the $b$-modes. But then the energy spectrum becomes arbitrarily negative by exciting more and more such states from the vacuum because of the minus sign in $H \sim-b^{\dagger} b$ !
We therefore conclude that this procedure results in

- either loss of unitarity due to appearance of negative-norm states
- or unboundness from below of $H$, i.e. instability of the vacuum.

The origin of this problem lies in the fact that the signs in the spinor theory conspire such that to establish a commutation relation of the schematic form $\left[\psi, \psi^{\dagger}\right] \sim 1$ we must impose

$$
\begin{equation*}
\left[a, a^{\dagger}\right] \sim 1 \text { and }\left[b^{\dagger}, b\right] \sim 1 \rightarrow\left[b, b^{\dagger}\right] \sim-1 . \tag{3.142}
\end{equation*}
$$

If instead we impose a relation symmetric in $\psi$ and $\psi^{\dagger}$ this minus sign for the $b$-mode relation would not occur. The task is therefore to promote the classical Poisson-bracket relations not to operator commutation relations, but to an analogous 'bracket' which is symmetric in both entries. The simplest such bracket is the anti-commutator. It turns out that this procedure is successful.

### 3.6.2 Using the anti-commutator

The correct procedure for quantisation of spin- $\frac{1}{2}$ fields is to impose the canonical anti-commutation relations

$$
\begin{align*}
& \left\{\psi^{A}(\vec{x}), \psi_{B}^{\dagger}\left(\vec{x}^{\prime}\right)\right\}=\delta_{B}^{A} \delta^{(3)}(\vec{x}-\vec{x}), \\
& \left\{\psi^{A}(\vec{x}), \psi^{B}\left(\vec{x}^{\prime}\right)\right\}=0=\left\{\psi_{A}^{\dagger}(\vec{x}), \psi_{B}^{\dagger}(\vec{x})\right\}, \tag{3.143}
\end{align*}
$$

where $\{A, B\}=A B+B A=\{B, A\}$. This induces the mode relations

$$
\begin{align*}
\left\{a_{r}(\vec{p}), a_{s}^{\dagger}(\vec{q})\right\} & =(2 \pi)^{3} \delta_{r s} \delta^{(3)}(\vec{p}-\vec{q}), \\
\left\{b_{r}(\vec{p}), b_{s}^{\dagger}(\vec{q})\right\} & =(2 \pi)^{3} \delta_{r s} \delta^{(3)}(\vec{p}-\vec{q}),  \tag{3.144}\\
\left\{a_{r}(\vec{p}), b_{s}^{(\dagger)}(\vec{q})\right\} & =0 .
\end{align*}
$$

Starting from $H \sim \sum a^{\dagger} a-b b^{\dagger}$ we now find

$$
\begin{equation*}
H=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} E_{p} \sum_{s}\left[a_{s}^{\dagger}(\vec{p}) a_{s}(\vec{p})+b_{s}^{\dagger}(\vec{p}) b_{s}(\vec{p})-(2 \pi)^{3} \delta(0)\right] . \tag{3.145}
\end{equation*}
$$

The divergent vacuum energy has opposite sign compared to a scalar theory. In theories with scalars and spin- $\frac{1}{2}$ fields cancellations in the vacuum energy are indeed possible. ${ }^{4}$

Since for this Hamiltonian the anti-commutation relations still imply the commutation relations (3.135) and (3.136), the vacuum is again defined by $a_{s}(\vec{p})|0\rangle=0=b_{s}(\vec{p})|0\rangle$. From this vacuum we define the Fock space of $a$ - and $b$-mode excitations. Let us start with the $a$-modes, which we will call the particle sector.

## 1-particle states

The state $|\vec{p}, s\rangle:=\sqrt{2 E_{p}} a_{s}^{\dagger}(\vec{p})|0\rangle$ is a 1-particle state with momentum $\vec{p}$, energy $E_{p}=\sqrt{\vec{p}^{2}+m^{2}}$ and spin $s$ in the $x_{3}$-direction, normalized such that

$$
\begin{equation*}
\langle\vec{p}, s \mid \vec{q}, r\rangle=2 E_{p}(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{q}) \delta^{r s} . \tag{3.146}
\end{equation*}
$$

## $N$-particle states

The state

$$
\begin{equation*}
\left|p_{1}, s_{1} ; \ldots ; p_{N}, s_{N}\right\rangle=\prod_{i=1}^{N} \sqrt{2 E_{p_{i}}} a_{s_{1}}^{\dagger}\left(\vec{p}_{1}\right) \ldots a_{s_{N}}^{\dagger}\left(\vec{p}_{N}\right)|0\rangle \tag{3.147}
\end{equation*}
$$

is an $N$-particle momentum eigenstate. This allows us, in complete analogy to the scalar field, to state the following theorem

The wavefunction of $N$-particle states of spin $\frac{1}{2}$ particles is anti-symmetric under particle exchange.

Indeed, if we exchange two particles, we pick up a minus sign due to the anti-commutation relations:

$$
\begin{equation*}
a_{s_{i}}^{\dagger}\left(\vec{p}_{i}\right) a_{s_{j}}^{\dagger}\left(\vec{p}_{j}\right)=-a_{s_{j}}^{\dagger}\left(\vec{p}_{j}\right) a_{s_{i}}^{\dagger}\left(\vec{p}_{i}\right) . \tag{3.148}
\end{equation*}
$$

This leads to the following Corollary

Spin $\frac{1}{2}$ particles obey Fermi-statistics, i.e. they are fermions.

In particular they obey the Pauli exclusion principle:

No two fermionic states of exactly the same quantum numbers are possible.

[^18]The Pauli exclusion principle is again a result of the anti-commutation relation because

$$
\begin{equation*}
a_{s}^{\dagger}(\vec{p}) a_{s}^{\dagger}(\vec{p})|0\rangle=0 . \tag{3.149}
\end{equation*}
$$

This exemplifies the much more general Spin-Statistics-Theorem ${ }^{5}$ :

Lorentz invariance, positivity of energy, unitarity and causality imply that:

- Particles of half-integer spin are fermions and
- particles of integer spin are bosons.

Excitations $b_{s}^{\dagger}(\vec{p})|0\rangle$ describe the corresponding anti-particles. Indeed the Lagrangian

$$
\begin{equation*}
\mathcal{L}=\bar{\psi}\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi \tag{3.150}
\end{equation*}
$$

enjoys a global $U(1)$ symmetry because it is invariant under

$$
\begin{equation*}
\psi \mapsto e^{-i \alpha} \psi, \quad \bar{\psi} \mapsto \bar{\psi} e^{i \alpha}, \quad \alpha \in \mathbb{R} \tag{3.151}
\end{equation*}
$$

The associated conserved Noether current will be found in the tutorial to take the form

$$
\begin{equation*}
j^{\mu}=\bar{\psi} \gamma^{\mu} \psi \tag{3.152}
\end{equation*}
$$

with Noether charge

$$
\begin{equation*}
Q=\int \mathrm{d}^{3} x j^{0}=\int \mathrm{d}^{3} x \psi^{\dagger} \psi=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \sum_{s}\left(a_{s}^{\dagger}(\vec{p}) a_{s}(\vec{p})-b_{s}^{\dagger}(\vec{p}) b_{s}(\vec{p})\right), \tag{3.153}
\end{equation*}
$$

after dropping a normal ordering constant. The charge acts on the 1-particle state as follows:

$$
\begin{align*}
& Q a_{s}^{\dagger}(\vec{p})|0\rangle=+a_{s}^{\dagger}(\vec{p})|0\rangle, \text { thus defining a fermion, }  \tag{3.154}\\
& Q b_{s}^{\dagger}(\vec{p})|0\rangle=-b_{s}^{\dagger}(\vec{p})|0\rangle, \text { thus defining an anti-fermion. }
\end{align*}
$$

Finally, a careful analysis of the angular momentum operator via Noether's theorem reveals that

$$
\begin{equation*}
J_{x_{3}} a_{s}^{\dagger}(\vec{p}=0)|0\rangle=s a_{s}^{\dagger}(\vec{p}=0)|0\rangle, \tag{3.155}
\end{equation*}
$$

but

$$
\begin{equation*}
J_{x_{3}} b_{s}^{\dagger}(\vec{p}=0)|0\rangle=-s b_{s}^{\dagger}(\vec{p}=0)|0\rangle, \tag{3.156}
\end{equation*}
$$

with $s= \pm \frac{1}{2}$, For details of the derivation see Peskin-Schröder, page 61. This shows that $b_{s}^{\dagger}(\vec{p}=0)|0\rangle$ has spin (in $x_{3}$-direction) $-s$, while $a_{s}^{\dagger}(\vec{p}=0)|0\rangle$ has spin $+s$.

[^19]
### 3.7 Propagators

As for the scalar fields, we now move to the Heisenberg picture by considering the time-dependent free fields (with free mass denoted by $m_{0}$ to avoid confusion)

$$
\begin{align*}
\psi(x) & =\sum_{s} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left[a_{s}(\vec{p}) u_{s}(\vec{p}) e^{-i p \cdot x}+b_{s}^{\dagger}(\vec{p}) v_{s}(\vec{p}) e^{i p \cdot x}\right]  \tag{3.157}\\
\psi^{\dagger}(x) & =\sum_{s} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left[b_{s}(\vec{p}) v_{s}^{\dagger}(\vec{p}) e^{-i p \cdot x}+a_{s}^{\dagger}(\vec{p}) u_{s}^{\dagger}(\vec{p}) e^{i p \cdot x}\right]
\end{align*}
$$

which, as noted already, satisfy the free Dirac equation as an operator equation.

To examine causality of the theory we define the anti-commutator

$$
\begin{equation*}
S_{B}^{A}(x-y):=\left\{\psi^{A}(x), \bar{\psi}_{B}(y)\right\} \tag{3.158}
\end{equation*}
$$

and compute

$$
\begin{align*}
S(x-y)= & \sum_{s, r} \int \frac{\mathrm{~d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{E_{p}}} \int \frac{\mathrm{~d}^{3} q}{(2 \pi)^{3}} \frac{1}{\sqrt{E_{q}}} \\
& \times[\underbrace{\left\{a_{s}(\vec{p}), a_{r}^{\dagger}(\vec{q})\right\}}_{\delta_{s r}(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{q})} u_{s}(\vec{p}) \bar{u}_{r}(\vec{q}) e^{-i p \cdot x} e^{i q \cdot y}  \tag{3.159}\\
& \left.+\left\{b_{s}^{\dagger}(\vec{p}), b_{r}(\vec{q})\right\} v_{s}(\vec{p}) \bar{v}_{r}(\vec{q}) e^{i p \cdot x} e^{-i q \cdot y}\right]
\end{align*}
$$

The identities

$$
\begin{equation*}
\sum_{s} u_{s}(\vec{p}) \bar{u}_{s}(\vec{p})=\gamma \cdot p+m_{0} \text { and } \sum_{s} v_{s}(\vec{p}) \bar{v}_{s}(\vec{p})=\gamma p-m_{0} \tag{3.160}
\end{equation*}
$$

imply that

$$
\begin{equation*}
S(x-y)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}\left[\left(\gamma \cdot p+m_{0}\right) e^{-i p \cdot(x-y)}+\left(\gamma \cdot p-m_{0}\right) e^{-i p \cdot(y-x)}\right] \tag{3.161}
\end{equation*}
$$

or in a more compact form

$$
\begin{equation*}
S(x-y)=\left(i \gamma^{\mu} \partial_{x^{\mu}}+m\right)\left[D^{(0)}(x-y)-D^{(0)}(y-x)\right] \tag{3.162}
\end{equation*}
$$

Here

$$
\begin{equation*}
D^{(0)}(x-y)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}} e^{-i p \cdot(x-y)} \tag{3.163}
\end{equation*}
$$

is the propagator from the scalar theory with mass $m_{0}$. In particular this implies that $S(x-y)=0$ for $(x-y)^{2}<0$. This in turn guarantees that $\left[O_{1}(x), O_{2}(y)\right]=0$ for $(x-y)^{2}<0$ for $O_{i}(x)$ any local expression of fermion bilinears $\bar{\psi} \psi$. Since all physical observables are bosonic this establishes causality of the Dirac theory.

The time-ordering symbol in the fermionic theory is defined as

$$
T(\psi(x) \bar{\psi}(y))=\left\{\begin{array}{c}
\psi(x) \bar{\psi}(y) \text { if } x^{0} \geq y^{0}  \tag{3.164}\\
-\bar{\psi}(y) \psi(x) \text { if } y^{0}>x^{0}
\end{array}\right.
$$

Note the crucial minus sign. It is required because if $(x-y)^{2}<0$, we have

$$
\begin{equation*}
\psi(x) \bar{\psi}(y)=-\bar{\psi}(y) \psi(x), \tag{3.165}
\end{equation*}
$$

because $S(x-y)=0$ for $(x-y)^{2}<0$. Now, for $(x-y)^{2}<0$ the question of whether $x^{0} \geq y^{0}$ or $x^{0}<y^{0}$ depends on the Lorentz frame we have chosen. To arrive at a Lorentz frame independent definition of the time-ordering symbol $T$ the expression for $T(\psi(x) \bar{\psi}(y))$ for $x^{0} \geq y^{0}$ and $y^{0} \geq x^{0}$ must agree.
The Feynman propagator is

$$
\begin{equation*}
S_{F}(x-y)=\langle 0| T \psi(x) \bar{\psi}(y)|0\rangle, \tag{3.166}
\end{equation*}
$$

while

$$
\begin{equation*}
\langle 0| T \psi(x) \psi(y)|0\rangle=0=\langle 0| T \bar{\psi}(x) \bar{\psi}(y)|0\rangle . \tag{3.167}
\end{equation*}
$$

Now by the usual tricks one evaluates

$$
\begin{align*}
S_{F}(x-y)= & \Theta\left(x^{0}-y^{0}\right)\langle 0| \psi(x) \bar{\psi}(y)|0\rangle-\Theta\left(y^{0}-x^{0}\right)\langle 0| \bar{\psi}(y) \psi(x)|0\rangle \\
= & \Theta\left(x^{0}-y^{0}\right) \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}\left(\gamma \cdot p+m_{0}\right) e^{-i p \cdot(x-y)} \\
& -\Theta\left(y^{0}-x^{0}\right) \int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{2 E_{p}}\left(\gamma \cdot p-m_{0}\right) e^{+i p \cdot(x-y)}  \tag{3.168}\\
= & \left(i \gamma \cdot \partial_{x}+m\right) \underbrace{D_{F}^{(0)}(x-y)}_{\text {free scalar theory }},
\end{align*}
$$

with

$$
\begin{equation*}
D_{F}^{(0)}(x-y)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i}{p^{2}-m_{0}^{2}+i \epsilon} e^{-i p \cdot(x-y)} . \tag{3.169}
\end{equation*}
$$

This gives

$$
\begin{equation*}
S_{F}(x-y)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i\left(\gamma \cdot p+m_{0}\right)}{p^{2}-m_{0}^{2}+i \epsilon} e^{-i p \cdot(x-y)} \tag{3.170}
\end{equation*}
$$

Due to $\left(\gamma \cdot p+m_{0}\right)\left(\gamma \cdot p-m_{0}\right)=p^{2}-m_{0}^{2}$ it is convenient to write

$$
\begin{equation*}
\frac{\left(\gamma \cdot p+m_{0}\right)}{p^{2}-m_{0}^{2}}=\left(\gamma \cdot p-m_{0}\right)^{-1} . \tag{3.171}
\end{equation*}
$$

The Feynman propagator $S_{F}(x-y)$ is a Green's function in that it represents one of the 4 possible solutions to

$$
\begin{equation*}
\left(i \gamma^{\mu} \partial_{\mu}-m\right) G(x-y)=i \delta^{(4)}(x-y) \tag{3.172}
\end{equation*}
$$

The interpretation and closure procedure in the complex plane for these are as in the scalar theory.

### 3.8 Wick's theorem and Feynman diagrams

The time ordering of several fields picks up a minus sign whenever 2 fermionic fields are exchanged, e.g.

$$
T\left(\psi_{1} \psi_{2} \psi_{3}\right)=\left\{\begin{array}{cc}
(-1) \psi_{1} \psi_{3} \psi_{2} & \text { if } x_{1}^{0}>x_{3}^{0}>x_{2}^{0}  \tag{3.173}\\
(-1)^{2} \psi_{3} \psi_{1} \psi_{2} & \text { if } x_{3}^{0}>x_{1}^{0}>x_{2}^{0} \\
(-1)^{3} \psi_{3} \psi_{2} \psi_{1} & \text { if } x_{3}^{0}>x_{2}^{0}>x_{1}^{0} \\
& \ldots
\end{array}\right.
$$

We define normal-ordered products as expressions with all creation operators to the left of all annihilation operators, where, unlike in the scalar theory, each exchange of two operators induces a minus sign, e.g.

$$
\begin{equation*}
: b_{s}(\vec{p}) a_{r}^{\dagger}(\vec{q}) b_{v}^{\dagger}(\vec{k}):=(-1)^{2} a_{r}^{\dagger}(\vec{q}) b_{v}^{\dagger}(\vec{k}) b_{s}(\vec{p}) \tag{3.174}
\end{equation*}
$$

and so on. Then

$$
\begin{equation*}
T(\psi(x) \bar{\psi}(y))=: \psi(x) \bar{\psi}(y):+\overleftarrow{\psi(x) \bar{\psi}(y)} \tag{3.175}
\end{equation*}
$$

with

$$
\begin{equation*}
\overline{\psi(x) \bar{\psi}(y)}=\langle 0| T \psi(x) \bar{\psi}(y)|0\rangle=S_{F}(x-y) \tag{3.176}
\end{equation*}
$$

and

$$
\begin{equation*}
\overline{\psi(x) \psi(y)}=0=\overline{\bar{\psi}(x) \bar{\psi}(y)} \tag{3.177}
\end{equation*}
$$

Direct computation confirms that Wick's theorem goes through, with the understanding that we include the minus signs from operator exchanges. For instance

$$
\begin{align*}
T\left(\psi_{1} \bar{\psi}_{2} \bar{\psi}_{3}\right)= & : \psi_{1} \bar{\psi}_{2} \bar{\psi}_{3}:+: \bar{\psi}_{1} \bar{\psi}_{2} \psi_{3}: \\
& +\underbrace{: \sqrt[\psi_{1} \bar{\psi}_{2} \bar{\psi}_{3}]{ }}_{-: \psi_{1} \bar{\psi}_{3} \bar{\psi}_{2}:=-S_{f}\left(x_{1}-x_{3}\right): \bar{\psi}_{2}:} . \tag{3.178}
\end{align*}
$$

With this in mind Wick's theorem becomes

$$
\begin{equation*}
T\left(\bar{\psi}_{1} \bar{\psi}_{2} \psi_{3} \ldots\right)=: \bar{\psi}_{1} \bar{\psi}_{2} \psi_{3} \ldots+\text { all contractions with signs : } \tag{3.179}
\end{equation*}
$$

In particular

$$
\begin{equation*}
\langle 0| T \prod_{i} \psi\left(x_{i}\right) \prod_{j} \bar{\psi}\left(\bar{x}_{j}\right)|0\rangle \neq 0 \tag{3.180}
\end{equation*}
$$

only for equal numbers of $\psi$ and $\bar{\psi}$ fields. Physically this just reflects charge conservation.

To compute a 2 n -point function of this type, we draw the corresponding Feynman diagrams, but now

- label the points $x_{i}$ associated with $\psi\left(x_{i}\right)$ and $\bar{x}_{j}$ associated with $\bar{\psi}\left(\bar{x}_{j}\right)$ separately,
- only connect $x_{i}$ with $\bar{x}_{j}$ and
- associate each directed line from $\bar{x}_{j}$ to $x_{i}$ with a propagator $S_{F}\left(x_{i}-\bar{x}_{j}\right)$.

Be sure to always draw the arrow from $\bar{x}_{j}$ to $x_{i}$ in order to account for the correct sign in $S_{F}\left(x_{i}-\bar{x}_{j}\right)$. Apart from an overall sign (which is typically unimportant because we will eventually take the square of the amplitude), the relative signs between the diagrams (which are important due to interference) equal the number of crossing lines. For instance, for

$$
\begin{equation*}
\langle 0| T \psi\left(x_{1}\right) \bar{\psi}\left(\bar{x}_{2}\right) \psi\left(x_{3}\right) \bar{\psi}\left(\bar{x}_{4}\right)|0\rangle \tag{3.181}
\end{equation*}
$$

this prescription gives (see Figure 3.1)

$$
\begin{equation*}
S_{F}\left(x_{1}-\bar{x}_{2}\right) S\left(x_{3}-\bar{x}_{4}\right)+(-1)^{1} S_{F}\left(x_{1}-\bar{x}_{4}\right) S\left(x_{3}-\bar{x}_{2}\right) . \tag{3.182}
\end{equation*}
$$



Figure 3.1: Possible Feynman diagrams

### 3.9 LSZ and Feynman rules

We will examine interacting spin $\frac{1}{2}$ fields in great detail in the context of Quantum Electrodynamics. Another example of an interesting interacting theory is Yukawa theory, which couples a spin $\frac{1}{2}$ field to a real boson via a cubic coupling. Its form is given in the tutorial.

In this section we only briefly summarise the logic behind the computation of scattering amplitudes with spin $\frac{1}{2}$. As in the scalar theory, in the presence of interactions we define asymptotic in- and out-fields satisfying the free Dirac equation with mass $m \neq m_{0}$, where $m_{0}$ is the mass in the free Dirac action. We then express the creation and annihilation modes by the in- and out-fields, e.g. for the
in-fields

$$
\begin{align*}
& a_{\mathrm{in}, s}(\vec{q})=\frac{1}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x \bar{u}_{s}(\vec{q}) e^{i q \cdot x} \gamma^{0} \psi_{\text {in }}(x), \\
& a_{\mathrm{in}, s}^{\dagger}(\vec{q})=\frac{1}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x \bar{\psi}_{\text {in }}(x) \gamma^{0} e^{-i q \cdot x} u_{s}(\vec{q}), \\
& b_{\text {in }, s}(\vec{q})=\frac{1}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x \bar{\psi}_{\text {in }}(x) \gamma^{0} e^{i q \cdot x} v_{s}(\vec{q}),  \tag{3.183}\\
& b_{\mathrm{in}, s}^{\dagger}(\vec{q})=\frac{1}{\sqrt{2 E_{q}}} \int \mathrm{~d}^{3} x \bar{v}_{s}(\vec{q}) e^{-i q \cdot x} \gamma^{0} \psi_{\text {in }}(x) .
\end{align*}
$$

Using these one can perform exactly the same LSZ reduction procedure as in the scalar field case to extract the $S$-matrix. In this process we make heavy use of the equations

$$
\begin{equation*}
(\gamma \cdot p-m) u(\vec{p})=0, \quad(\gamma \cdot p+m) v(\vec{p})=0 \tag{3.184}
\end{equation*}
$$

where now $m$ is the fully renormalized physical mass. Consider incoming fermions $|q, s,+\rangle$ and antifermions $\left|q^{\prime} s^{\prime},-\right\rangle$ and outgoing fermions $\langle p, r,+|$ and anti-fermions $\left\langle q^{\prime}, r^{\prime},-\right|$. The final result for the $S$-matrix element is

$$
\begin{align*}
& \left.\left\langle\ldots(p, r,+) \ldots\left(p^{\prime} r^{\prime},-\right) \ldots\right| S\left|\ldots(q, s,+) \ldots\left(q^{\prime}, s^{\prime},-\right) \ldots\right\rangle\right|_{\text {connected }} \\
= & \left((-i Z)^{-\frac{1}{2}}\right)^{n}\left((i Z)^{-\frac{1}{2}}\right)^{n^{\prime}} \int \mathrm{d}^{4} x \ldots \int \mathrm{~d}^{4} x^{\prime} \ldots \int \mathrm{d}^{4} y \ldots \int \mathrm{~d}^{4} y^{\prime}  \tag{3.185}\\
& \times e^{-i\left(q \cdot x+q^{\prime} \cdot x^{\prime}-p \cdot y-p^{\prime} \cdot y^{\prime}+\ldots\right)} \bar{u}_{r}(\vec{p})\left(i \gamma \cdot \partial_{y}-m\right) \ldots \bar{v}_{s^{\prime}}\left(\vec{q}^{\prime}\right)\left(i \gamma \cdot \partial_{x^{\prime}}-m\right) \ldots \\
& \times\langle\Omega| T \ldots \bar{\psi}\left(y^{\prime}\right) \ldots \psi(y) \ldots \bar{\psi}(x) \ldots \psi\left(x^{\prime}\right) \ldots|\Omega\rangle\left(-i \gamma \cdot \overleftarrow{\partial_{x}}-m\right) u_{s}(\vec{q}) \ldots\left(-i \gamma \cdot \overleftarrow{\partial_{y}}-m\right) v_{r^{\prime}}\left(\vec{p}^{\prime}\right) .
\end{align*}
$$

Thus to compute the S-matrix we compute the Fourier transform of the amputated fully connected associated Feynman diagram, where for each external particle we include

- $u_{s}(\vec{q})$ for an incoming particle of spin $s$,
- $\bar{v}_{s^{\prime}}\left(\vec{q}^{\prime}\right)$ for an incoming anti-particle of spin $-s^{\prime}$,
- $\bar{u}_{r}(\vec{p})$ for an outgoing particle of $\operatorname{spin} r$,
- $v_{r^{\prime}}\left(\vec{p}^{\prime}\right)$ for an outgoing anti-particle of spin $-r^{\prime}$.

Alternatively, the appearance of these spinor polarisations can also be deduced by the Interaction picture procedure discussed in the tutorials. We will exemplify this for Yukawa theory in the tutorial and later in the course for Quantum Electrodynamics.

## Chapter 4

## Quantising spin 1-fields

### 4.1 Classical Maxwell-theory

Let us first recall the main aspects of classical Maxwell theory.

- The classical Maxwell equations are

$$
\begin{array}{ll}
\vec{\nabla} \cdot \vec{E}=\rho, & \vec{\nabla} \times \vec{E}=-\frac{\partial \vec{B}}{\partial t}  \tag{4.1}\\
\vec{\nabla} \times \vec{B}=\vec{j}+\frac{\partial \vec{E}}{\partial t}, & \vec{\nabla} \cdot \vec{B}=0,
\end{array}
$$

where the sources are subject to local charge conservation, i.e.

$$
\begin{equation*}
\frac{\partial}{\partial t} \rho+\vec{\nabla} \cdot \vec{j}=0 \tag{4.2}
\end{equation*}
$$

- By virtue of the inhomogeneous Maxwell equations and Helmholtz's theorem the fields $\vec{E}$ and $\vec{B}$ can locally be expressed as

$$
\begin{equation*}
\vec{E}=-\vec{\nabla} \phi-\frac{\partial \vec{A}}{\partial t}, \quad \vec{B}=\vec{\nabla} \times \vec{A} . \tag{4.3}
\end{equation*}
$$

- This description in terms of the scalar potential $\phi$ and the vector potential $\vec{A}$ is redundant because $\vec{E}$ and $\vec{B}$ are invariant under a gauge transformation, i.e. a transformation

$$
\begin{align*}
\phi(x) & \rightarrow \phi(x)+\frac{\partial}{\partial t} \alpha(x)  \tag{4.4}\\
\vec{A}(x) & \rightarrow \vec{A}(x)-\vec{\nabla} \alpha(x)
\end{align*}
$$

- To establish a Lorentz invariant formulation we introduce the 4-vector gauge potential $A^{\mu}$ and the 4-current $j^{\mu}$ as

$$
\begin{equation*}
A^{\mu}=\binom{\phi}{\vec{A}}, \quad j^{\mu}=\binom{\rho}{\vec{j}} . \tag{4.5}
\end{equation*}
$$

The fields $\vec{E}$ and $\vec{B}$ are really components of the field strength tensor $F^{\mu \nu}$,

$$
\begin{equation*}
F^{\mu v}=\partial^{\mu} A^{\nu}-\partial^{\nu} A^{\mu} \tag{4.6}
\end{equation*}
$$

which in matrix notation reads

$$
F_{\mu \nu}=\left(\begin{array}{cccc}
0 & E_{1} & E_{2} & E_{3}  \tag{4.7}\\
-E_{1} & 0 & -B_{3} & B_{2} \\
-E_{2} & B_{3} & 0 & -B_{1} \\
-E_{3} & -B_{2} & B_{1} & 0
\end{array}\right) .
$$

- Using the field strength tensor the inhomogenous Maxwell equations can be written as

$$
\begin{equation*}
\partial_{\mu} F^{\mu \nu}=j^{\nu} . \tag{4.8}
\end{equation*}
$$

The homogenous Maxwell equations are automatically satisfied because they were used to express $\vec{E}$ and $\vec{B}$ in terms of the potentials. Indeed they correspond to the Bianchi identity

$$
\begin{equation*}
\partial_{[\mu} F_{\nu \rho]}=0, \tag{4.9}
\end{equation*}
$$

where [] denotes all cyclic permutations. Note that $\partial_{\mu} j^{\mu}=0$ follows as a consistency condition because $\partial_{\nu} \partial_{\mu} F^{\mu \nu}=0$ since we contract the symmetric tensor $\partial_{\nu} \partial_{\mu}$ with the anti-symmetric tensor $F^{\mu \nu}$.

- We stress again that $F^{\mu v}$ and thus $\vec{E}$ and $\vec{B}$ are invariant under a local gauge transformation

$$
\begin{equation*}
A^{\mu}(x) \rightarrow A^{\mu}(x)+\partial^{\mu} \alpha(x) . \tag{4.10}
\end{equation*}
$$

Configurations related by gauge transformations are physically equivalent. Gauge symmetries merely denote a redundancy in the description of the system. To determine the true local physical degrees of freedom we must be sure to divide out by this redundancy. This will be the main difficulty in quantising the system.

- The Maxwell equations follow as the equation of motion of $A^{\mu}$ from the action

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left(-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}-A_{\mu}(x) j^{\mu}(x)\right) . \tag{4.11}
\end{equation*}
$$

Note that $S$ is gauge invariant if and only if $\partial_{\mu} j^{\mu}=0$. The equation of motion can be rewritten as

$$
\begin{equation*}
\square A^{\mu}-\partial^{\mu}\left(\partial_{\nu} A^{\nu}\right)=j^{\mu} \tag{4.12}
\end{equation*}
$$

- Due to gauge invariance we can always pick $A^{\mu}$ such that

$$
\begin{equation*}
\partial_{\mu} A^{\mu}=0 . \tag{4.13}
\end{equation*}
$$

This partially fixes the gauge in Lorenz gauge, but we are still free to perform a residual gauge transformation

$$
\begin{equation*}
A^{\mu} \rightarrow A^{\mu}+\partial^{\mu} \phi \quad \text { with } \quad \square \phi=0 \tag{4.14}
\end{equation*}
$$

without violating the Lorenz gauge condition. In Lorenz gauge the equation of motion is

$$
\begin{equation*}
\square A^{\mu}=j^{\mu} . \tag{4.15}
\end{equation*}
$$

- Lorenz gauge can be implemented by adding a Lagrange multiplier term in the action:

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left(-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}-\frac{\lambda}{2}(\partial \cdot A)^{2}-A_{\mu} j^{\mu}\right) \tag{4.16}
\end{equation*}
$$

Therefore we now have two equations of motion, namely for $A^{\mu}$

$$
\begin{equation*}
\square A^{\mu}-(1-\lambda) \partial^{\mu}(\partial \cdot A)=j^{\mu} \tag{4.17}
\end{equation*}
$$

and for $\lambda$

$$
\begin{equation*}
\partial \cdot A=0, \tag{4.18}
\end{equation*}
$$

where by equation of motion for $\lambda$ we mean that the variation of $S$ with respect to $\lambda$ is proportional to $\partial \cdot A$.

### 4.2 Canonical quantisation of the free field

We now set the current $j^{\mu}=0$ and consider the free gauge potential. The free non-gauge fixed Lagrangian is

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}=-\frac{1}{4}\left(\partial_{\mu} A_{\nu}-\partial_{\nu} A_{\mu}\right)\left(\partial^{\mu} A^{\nu}-\partial^{\nu} A^{\mu}\right) \tag{4.19}
\end{equation*}
$$

It is not suitable for quantisation because the momentum density canonically conjugate to $A^{\mu}$ is

$$
\begin{equation*}
\Pi_{\mu}=\frac{\partial \mathcal{L}}{\partial \dot{\mathcal{A}}^{\mu}}=F_{\mu 0}, \tag{4.20}
\end{equation*}
$$

and since

$$
\begin{equation*}
\Pi_{0}=0, \tag{4.21}
\end{equation*}
$$

$\left(A^{\mu}, \Pi_{\mu}\right)$ are no good canonical variables. Instead quantisation starts from the gauge fixed Lagrangian

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}-\frac{\lambda}{2}(\partial \cdot A)^{2}, \tag{4.22}
\end{equation*}
$$

with

$$
\begin{equation*}
\Pi_{\mu}=F_{\mu 0}-\lambda \eta_{\mu 0}(\partial \cdot A) . \tag{4.23}
\end{equation*}
$$

We could proceed for a general Lagrange multiplier $\lambda$, but for simplicity we set

$$
\begin{equation*}
\lambda=1 \text { corresponding to Feynman gauge. } \tag{4.24}
\end{equation*}
$$

Explicitly the Lagrangian in Feynman gauge is

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{2} \partial_{\mu} A_{\nu} \partial^{\mu} A^{\nu} \quad \text { together with } \partial \cdot A=0 \tag{4.25}
\end{equation*}
$$

Since the Lagrangian multiplier $\lambda$ has been integrated out by setting $\lambda=1$, its equation of motion $\partial \cdot A=0$ must now be imposed by hand as a constraint.

Note two important facts about the gauge fixed Lagrangian:

- The Langrangian $\mathcal{L}=-\frac{1}{2} \partial_{\mu} A_{\nu} \partial^{\mu} A^{\nu}$ is simply the Lagrangian of 3 free massless scalars $A^{i}$, $\mathrm{i}=1,2,3$, but for $\mu=0$ the sign of the kinetic terms is wrong. This will be important in the sequel.
- The extra constraint $\partial \cdot A=0$ is a consequence of the underlying gauge symmetry of the system and will ensure that a consistent quantization is possible despite the wrong sign for $\mu=0$.

The equation of motion for $A^{\mu}$ which follows from (4.25) is

$$
\begin{equation*}
\square A^{\mu}=0 \quad \text { together with } \partial \cdot A=0 . \tag{4.26}
\end{equation*}
$$

With (4.25) as a starting point, the canonical momentum density is

$$
\begin{equation*}
\Pi_{\mu}=-\dot{A}_{\mu} \tag{4.27}
\end{equation*}
$$

We quantise the system by promoting $A^{\mu}$ and $\Pi_{v}$ to Heisenberg picture fields with canonical equaltime commutators

$$
\begin{equation*}
\left[A^{\mu}(t, \vec{x}), \Pi_{\nu}(t, \vec{y})\right]=i \delta^{\mu}{ }_{\nu} \delta^{(3)}(\vec{x}-\vec{y}) \tag{4.28}
\end{equation*}
$$

and therefore

$$
\begin{align*}
& {\left[A^{\mu}(t, \vec{x}), \dot{A}^{\nu}(t, \vec{y})\right]=-i \eta^{\mu \nu} \delta^{(3)}(\vec{x}-\vec{y}),}  \tag{4.29}\\
& {\left[A^{\mu}(t, \vec{x}), A^{v}(t, \vec{y})\right]=0=\left[\dot{A}^{\mu}(t, \vec{x}), \dot{A}^{v}(t, \vec{y})\right] .}
\end{align*}
$$

As observed above there is an odd minus sign for $\mu=v=0$. Despite this issue we proceed and consider the mode expansion

$$
\begin{equation*}
A^{\mu}(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}} \sum_{\lambda=0}^{3} \epsilon^{\mu}(\vec{p}, \lambda)\left[a_{\lambda}(\vec{p}) e^{-i p \cdot x}+a_{\lambda}^{\dagger}(\vec{p}) e^{i p \cdot x}\right], \tag{4.30}
\end{equation*}
$$

which is compatible with the equation of motion $\square A^{\mu}=0$ as an operator equation provided

$$
\begin{equation*}
p^{2}=0 \Rightarrow E_{p}=|\vec{p}| . \tag{4.31}
\end{equation*}
$$

Furthermore the vectors $\epsilon^{\mu}(\vec{p}, \lambda), \lambda=0,1,2,3$ are the 4 linearly independent real polarisation vectors whose definition depends on the value of the lightlike vector $p^{\mu}$ (satisfying $p^{2}=0$ ). Our conventions to define these are as follows: Let $n^{\mu}$ denote the time axis such that $n^{2}=1$.

- $\epsilon^{\mu}(\vec{p}, 0) \equiv n^{\mu}$ is the timelike or scalar polarisation vector.
- $\epsilon^{\mu}(\vec{p}, i)$ for $i=1,2$ are called transverse polarisation vectors. They are defined via $\epsilon(\vec{p}, i) \cdot n=$ $0=\epsilon(\vec{p}, i) \cdot p$ and

$$
\begin{equation*}
\epsilon(\vec{p}, i) \cdot \epsilon(\vec{p}, j)=-\delta_{i j} . \tag{4.32}
\end{equation*}
$$

- $\epsilon^{\mu}(\vec{p}, 3)$ is called longitudinal polarisation and is defined via $\epsilon(\vec{p}, 3) \cdot n=0=\epsilon(\vec{p}, 3) \cdot \epsilon(\vec{p}, i)$ for $i=1,2$ and $\epsilon(\vec{p}, 3)^{2}=-1$. Since $p^{2}=0$ it can therefore be expressed as

$$
\begin{equation*}
\epsilon(\vec{p}, 3)=\frac{p-n(p \cdot n)}{p+n} . \tag{4.33}
\end{equation*}
$$

So altogether

$$
\begin{equation*}
\epsilon(\vec{p}, \lambda) \cdot \epsilon\left(\vec{p}, \lambda^{\prime}\right)=\eta_{\lambda, \lambda^{\prime}} . \tag{4.34}
\end{equation*}
$$

We stress that this basis of polarisation vectors depends on the concrete momentum vector $p$ with $p^{2}=0$. Consider e.g. a momentum vector $p^{\mu}=(1,0,0,1)^{T}$, then

$$
\epsilon(\vec{p}, 0)=\left(\begin{array}{l}
1  \tag{4.35}\\
0 \\
0 \\
0
\end{array}\right), \quad \epsilon^{\mu}(\vec{p}, 1)=\left(\begin{array}{l}
0 \\
1 \\
0 \\
0
\end{array}\right), \quad \epsilon^{\mu}(\vec{p}, 2)=\left(\begin{array}{l}
0 \\
0 \\
1 \\
0
\end{array}\right), \quad \epsilon^{\mu}(\vec{p}, 3)=\left(\begin{array}{l}
0 \\
0 \\
0 \\
1
\end{array}\right) .
$$

The canonical commutation relations imply for the modes

$$
\begin{equation*}
\left[a_{\lambda}(\vec{p}), a_{\lambda^{\prime}}^{\dagger}\left(\vec{p}^{\prime}\right)\right]=-\eta_{\lambda \lambda^{\prime}}(2 \pi)^{3} \delta^{(3)}\left(\vec{p}-\vec{p}^{\prime}\right) \tag{4.36}
\end{equation*}
$$

and

$$
\begin{equation*}
\left[a_{\lambda}(\vec{p}), a_{\lambda^{\prime}}\left(\vec{p}^{\prime}\right)\right]=0=\left[a_{\lambda}^{\dagger}(\vec{p}), a_{\lambda^{\prime}}^{\dagger}\left(\vec{p}^{\prime}\right)\right] . \tag{4.37}
\end{equation*}
$$

Note again the minus sign for timelike modes $\lambda=\lambda^{\prime}=0$.
The Hamiltonian is

$$
\begin{equation*}
H=\int \mathrm{d}^{3} x\left(-\dot{A}^{\mu} \dot{A}_{\mu}-\mathcal{L}\right)=\frac{1}{2} \int \mathrm{~d}^{3} x\left[-\dot{A}^{\mu} \dot{A}_{\mu}+\partial_{i} A_{\mu} \partial^{i} A^{\mu}\right] \tag{4.38}
\end{equation*}
$$

and in modes

$$
\begin{align*}
H & =-\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}}|\vec{p}| \sum_{\lambda} a_{\lambda}^{\dagger}(\vec{p}) a_{\lambda}(\vec{p}) \epsilon^{\mu}(\vec{p}, \lambda) \epsilon^{\nu}(\vec{p}, \lambda) \eta_{\mu v}  \tag{4.39}\\
& =\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}}|\vec{p}|\left[\sum_{i} a_{i}^{\dagger}(\vec{p}) a_{i}(\vec{p})-a_{0}^{\dagger}(\vec{p}) a_{0}(\vec{p})\right]
\end{align*}
$$

after dropping the vacuum energy. This leads to the following commutation relations (valid $\forall \lambda$ )

$$
\begin{align*}
{\left[H, a_{\lambda}^{\dagger}(\vec{p})\right] } & =+|\vec{p}| a_{\lambda}^{\dagger}(\vec{p}), \\
{\left[H, a_{\lambda}(\vec{p})\right] } & =-|\vec{p}| a_{\lambda}(\vec{p}) . \tag{4.40}
\end{align*}
$$

We define again the vacuum $|0\rangle$ such that $a_{\lambda}(\vec{p})|0\rangle=0$ and the 1-particle states

$$
\begin{equation*}
|\vec{p}, \lambda\rangle:=\sqrt{2 E_{p}} a_{\lambda}^{\dagger}(\vec{p})|0\rangle \tag{4.41}
\end{equation*}
$$

as the states of momentum $\vec{p}$ and polarisation $\lambda$. The corresponding particles are called photons.

Two crucial problems remain though:

- The previous analysis seems to suggest that the theory gives rise to 4 independent degrees of freedom per momentum eigenstate, but from classical electrodynamics we only expect 2 transverse degrees of freedom.
- Timelike polarisation states have negative norm,

$$
\begin{equation*}
\langle\vec{p}, 0 \mid \vec{q}, 0\rangle \propto\langle 0|\left[a(\vec{p}, 0), a^{\dagger}(\vec{q}, 0)\right]|0\rangle=-(2 \pi)^{3} \delta^{(3)}(\vec{p}-\vec{q}) . \tag{4.42}
\end{equation*}
$$

Such negative norm states are called ghosts and spoil unitarity.

### 4.3 Gupta-Bleuler quantisation

The two above problems arose because the quantisation procedure so far is incomplete: The point is that the constraint $\partial \cdot A=0$ has not been implemented yet. It is impossible to implement this constraint as an operator equation for Heisenberg fields because then we would conclude

$$
\begin{equation*}
0 \stackrel{!}{=}\left[\partial_{\mu} A^{\mu}(t, \vec{x}), A^{v}(t, \vec{y})\right]=\left[\dot{A}^{0}(t, \vec{x}), A^{v}(t, \vec{y})\right]=i \eta^{0 v} \delta^{(3)}(\vec{x}-\vec{y}) . \tag{4.43}
\end{equation*}
$$

The idea of the Gupta-Bleuler formalism is to implement $\partial \cdot A=0$ not at the level of operators, but directly on the Hilbert space, i.e. as a defining constraint on the so-called physical states. The naive guess would be to require

$$
\begin{equation*}
\partial \cdot A|\phi\rangle \stackrel{!}{=} 0 \tag{4.44}
\end{equation*}
$$

for $|\phi\rangle$ to be a physical state. Let us decompose

$$
\begin{equation*}
A(x)=A^{+}(x)+A^{-}(x), \tag{4.45}
\end{equation*}
$$

with

$$
\begin{align*}
& A^{+}(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2|\vec{p}|}} \sum_{\lambda} \epsilon(\vec{p}, \lambda) a_{\lambda}(\vec{p}) e^{-i p \cdot x},  \tag{4.46}\\
& A^{-}(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2|\vec{p}|}} \sum_{\lambda} \epsilon(\vec{p}, \lambda) a_{\lambda}^{\dagger}(\vec{p}) e^{i p \cdot x}
\end{align*}
$$

so that the physical state condition would be

$$
\begin{equation*}
\left(\partial \cdot A^{+}+\partial \cdot A^{-}\right)|\phi\rangle \stackrel{!}{=} 0 . \tag{4.47}
\end{equation*}
$$

But this is still too strong because then not even the vacuum $|0\rangle$ would be such a physical state - after all

$$
\begin{equation*}
\partial \cdot A^{+}|0\rangle=0, \text { but } \partial \cdot A^{-}|0\rangle \neq 0 \tag{4.48}
\end{equation*}
$$

However, the milder constraint

$$
\begin{equation*}
\partial \cdot A^{+}|\phi\rangle \stackrel{!}{=} 0 \stackrel{!}{=}\langle\phi| \partial \cdot A^{-} \tag{4.49}
\end{equation*}
$$

suffices to guarantee

$$
\begin{equation*}
\langle\phi| \partial \cdot A|\phi\rangle=\langle\phi|\left(\partial \cdot A^{+}|\phi\rangle\right)+\left(\langle\phi| \partial \cdot A^{-}\right)|\phi\rangle=0 \tag{4.50}
\end{equation*}
$$

So in the spirit of Ehrenfest's theorem the classical relation $\partial \cdot A=0$ is realised as a statement about the expectation value $\langle\phi| \partial \cdot A|\phi\rangle=0$ in the quantum theory.
To summarise: Out of the naive Fock space we define the physical Hilbert space by

$$
\begin{equation*}
\phi \in \mathcal{H}_{\text {phys }} \leftrightarrow \partial \cdot A^{+}|\phi\rangle=0 \tag{4.51}
\end{equation*}
$$

This is the Gupta-Bleuler condition.

It suffices to construct the physical 1-particle states of definite momentum since $\mathcal{H}_{\text {phys }}$ is spanned by the tensor product of these. Consider a state $|\vec{p}, \zeta\rangle$, i.e. 1 photon of polarisation $\zeta^{\mu}$, where we define a general polarisation 4-vector

$$
\begin{equation*}
\zeta^{\mu}=\sum_{\lambda, \lambda^{\prime}} \alpha_{\lambda} \eta_{\lambda \lambda^{\prime}} \epsilon^{\mu}(\vec{p}, \lambda) \tag{4.52}
\end{equation*}
$$

The state $|\vec{p}, \zeta\rangle$ is defined as

$$
\begin{equation*}
|\vec{p}, \zeta\rangle:=\sqrt{2|\vec{p}|} \sum_{\lambda} \alpha_{\lambda} a_{\lambda}^{\dagger}(\vec{p})|0\rangle \tag{4.53}
\end{equation*}
$$

Therefore

$$
\begin{equation*}
\left\langle\vec{q}, \zeta \mid \vec{p}, \zeta^{\prime}\right\rangle=-(2 \pi)^{3} \cdot 2|\vec{p}| \delta^{(3)}(\vec{p}-\vec{q}) \zeta \cdot \zeta^{\prime} \tag{4.54}
\end{equation*}
$$

The physical state condition $\partial_{\mu} A^{+\mu}|\vec{p}, \zeta\rangle \stackrel{!}{=} 0$ implies that $p_{\mu} \zeta^{\mu}=0$ because

$$
\begin{align*}
\partial_{\mu} A^{+\mu}|\vec{p}, \zeta\rangle & =\int \frac{\mathrm{d}^{3} q}{(2 \pi)^{3}} \frac{\sqrt{2|\vec{p}|}}{\sqrt{2|\vec{q}|}} \sum_{\lambda}\left(-i q_{\mu}\right) \epsilon^{\mu}(\vec{q}, \lambda) e^{-i q \cdot x} \underbrace{a_{\lambda}(\vec{q}) \sum_{\gamma} \alpha_{\gamma} a_{\gamma}^{\dagger}(\vec{p})|0\rangle}_{=-(2 \pi)^{3} \sum_{\gamma} \alpha_{\gamma} \eta_{l \gamma} \delta^{(3)}(\vec{p}-\vec{q})|0\rangle} \\
& =i p_{\mu} \zeta^{\mu}|0\rangle \tag{4.55}
\end{align*}
$$

Therefore

$$
\begin{equation*}
|\vec{p}, \zeta\rangle \in \mathcal{H}_{\mathrm{phys}} \leftrightarrow p^{\mu} \zeta_{\mu}=0 \tag{4.56}
\end{equation*}
$$

Since $p^{2}=0$ for a massless photon, such $\zeta^{\mu}$ can be decomposed as

$$
\begin{equation*}
\zeta^{\mu}=\zeta_{T}^{\mu}+\zeta_{S}^{\mu} \tag{4.57}
\end{equation*}
$$

with $\zeta_{S}=c \cdot p, \zeta_{S}^{2}=0$ and $\vec{\zeta}_{T} \cdot \vec{p}_{T}=0, \zeta_{T}^{2}\langle 0,$. So $\mid \vec{p}, \zeta\rangle \in \mathcal{H}_{\text {phys }}$ can be written as

$$
\begin{equation*}
|\vec{p}, \zeta\rangle=\left|\vec{p}, \zeta_{T}\right\rangle+\left|\vec{p}, \zeta_{s}\right\rangle \tag{4.58}
\end{equation*}
$$

where

- $\left|\vec{p}, \zeta_{T}\right\rangle$ describes 2 transverse degrees of freedom of positive norm, $\|\left|\vec{p}, \zeta_{T}\right\rangle \|>0$,
- $\left|\vec{p}, \zeta_{S}\right\rangle$ describes 1 combined timelike and longitudinal degree of freedom of zero norm, $\|\left|\vec{p}, \zeta_{S}\right\rangle \|=$ 0.

For example consider

$$
p^{\mu}=\left(\begin{array}{l}
1  \tag{4.59}\\
0 \\
0 \\
1
\end{array}\right) \Rightarrow \zeta_{T}=\left(\begin{array}{c}
0 \\
\zeta^{1} \\
\zeta^{2} \\
0
\end{array}\right), \quad \zeta_{s}=c\left(\begin{array}{l}
1 \\
0 \\
0 \\
1
\end{array}\right)
$$

and $\left|\vec{p}, \zeta_{S}\right\rangle \propto\left(a_{0}^{\dagger}(\vec{p})-a_{3}^{\dagger}(\vec{p})\right)|0\rangle$. Up to now, the Gupta-Bleuler procedure has eliminated the negative norm states and left us with 3 polarisation states, but in fact one can prove the following theorem:

The state $\left|\vec{p}, \zeta_{S}\right\rangle$ decouples from all physical processes.

Such a zero-norm state that decouples from all physical processes is called spurious, hence the subscript $S$. The meaning of this decoupling of null states is as follows:

- For the free theory decoupling means that

$$
\begin{equation*}
\left\langle\vec{p}, \zeta_{S}\right| O\left|\vec{p}, \zeta_{S}\right\rangle=0 \tag{4.60}
\end{equation*}
$$

for all observables $O$. As an example it is easy to check that

$$
\begin{equation*}
\left\langle\vec{p}, \zeta_{S}\right| H\left|\vec{p}, \zeta_{S}\right\rangle=0 . \tag{4.61}
\end{equation*}
$$

Without loss of generality we take $p^{\mu}=(1,0,0,1)^{T}$ and $\left|\vec{p}, \zeta_{S}\right\rangle \propto\left(a_{0}^{\dagger}(\vec{p})-a_{3}^{\dagger}(\vec{p})\right)|0\rangle$ and confirm $\left\langle\vec{p}, \zeta_{S}\right| H\left|\vec{p}, \zeta_{S}\right\rangle=0$ by noting the structure of the Hamiltonian

$$
\begin{equation*}
H \sim \sum_{i} a_{i}^{\dagger} a_{i}-a_{0}^{\dagger} a_{0} \tag{4.62}
\end{equation*}
$$

and the relative minus sign in the commutation relations for timelike and spacelike modes.
Furthermore the spurious states decouple in the sense that $\left|\vec{p}, \zeta_{S}\right\rangle$ has zero overlap with $\left|\vec{p}, \zeta_{T}\right\rangle$ because $\zeta_{S} \cdot \zeta_{T}=0$.

- The decoupling statement becomes actually non-trivial in the presence of interactions: As long as the interactions respect gauge invariance, a spurious state $\left|\vec{p}, \zeta_{S}\right\rangle$ decouples from the $S$-matrix as an external (in or out) state. This follows from the Ward identities as will be discussed later.

The conclusion is that only the 2 transverse polarisations are physically relevant as external states. These have positive norm.
This does not mean that spurious states play no role at all:

- Spurious states do appear as internal states in $S$-matrix processes and are important for consistency of the amplitudes.
- Spurious states are required to establish Lorentz-invariance of the theory because $\zeta_{T}$ may pick up a spurious component $\zeta_{s}$ by going to a different Lorentz frame.

Let $\left|\psi_{S}\right\rangle$ denote any multi-photon state constructed entirely out of spurious photons with polarization $\zeta_{s}$. Since it decouples in the above sense, we can add and substract it without affecting any physical properties of a state. This establishes an equivalence relation on $\mathcal{H}_{\text {phys }}$ :

$$
\begin{equation*}
\left|\phi_{1}\right\rangle \sim\left|\phi_{2}\right\rangle \text { if } \exists\left|\psi_{s}\right\rangle:\left|\phi_{1}\right\rangle=\left|\phi_{2}\right\rangle+\left|\psi_{S}\right\rangle . \tag{4.63}
\end{equation*}
$$

This is the analogue of the residual gauge symmetry in classical theory: Indeed, one can show that ${ }^{1}$

$$
\begin{equation*}
\left\langle\psi_{s}\right| A_{\mu}(x)\left|\psi_{s}\right\rangle=\partial_{\mu} \Lambda(x) \tag{4.64}
\end{equation*}
$$

for a function $\Lambda(x)$ with

$$
\begin{equation*}
\square \Lambda(x)=0 . \tag{4.65}
\end{equation*}
$$

Therefore:

Adding $\left|\psi_{s}\right\rangle$ in $|\phi\rangle \rightarrow|\phi\rangle+\left|\psi_{s}\right\rangle$ is the quantum version of the residual transformation

$$
\begin{equation*}
A_{\mu} \rightarrow A_{\mu}+\partial_{\mu} \Lambda \text { with } \square \Lambda=0 \tag{4.66}
\end{equation*}
$$

in Lorenz gauge.

## A more general perspective on massless Spin 1 fields

Our starting point was not an arbitrary massless vector field $A^{\mu}(x)$, but the very specific gauge potential that arises in Maxwell-theory. More generally we might ask: Given a general massless vector field $A^{\mu}(x)$, how can we quantise it?

The definitive treatment of this question can be found in Weinberg, Vol. I, Chapter 8.1, which we urgently recommend. Following this reference, the arguments are:

- By Lorentz invariance alone, any massless vector field $A_{\mu}(x)$ must describe precisely two helicity or polarization states (the two transverse degrees of freedom we found above).

[^20]- On general grounds one can show that Lorentz vector fields describing two polarization states transform under a Lorentz transformation as

$$
\begin{equation*}
A^{\mu}(x) \rightarrow \Lambda_{v}^{\mu} A^{v}\left(\Lambda^{-1} x\right)+\partial^{\mu} \epsilon(x, \Lambda) \tag{4.67}
\end{equation*}
$$

for a spacetime-dependent function $\epsilon(x, \Lambda)$. Therefore Lorentz invariance requires invariance of the action under gauge transformations.

- The Lagrangian $\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}$ is the unique Lorentz invariant and gauge invariant Lagrangian for a massless free vector field (i.e. up to quadratic order).

This proves the general statement:
Massless vector field theories must be gauge theories.

### 4.4 Massive vector fields

Theories of massive vector bosons on the other hand are consistent despite the lack of gauge invariance: Consider the Lagrangian for "massive electrodynamics"

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}+\frac{1}{2} \mu^{2} A_{\mu} A^{\mu} \tag{4.68}
\end{equation*}
$$

with the classical equation of motion, the so-called Proca equation,

$$
\begin{equation*}
\partial_{\mu} F^{\mu \nu}+\mu^{2} A^{\nu}=0 . \tag{4.69}
\end{equation*}
$$

One can show that this Lagrangian is the unique Lorentz invariant Lagrangian for a free massive spin1 field (without any spin- 0 components - see Weinberg I, 7.5 for a proof). The mass term explicitly breaks gauge invariance. So it is not possible to arrange for $\partial \cdot A=0$ by gauge-fixing. However, the Proca equation implies

$$
\begin{equation*}
0=\underbrace{\partial_{\mu} \partial_{v} F^{\mu \nu}}_{=0}+\mu^{2} \partial_{v} A^{\nu} \Rightarrow \partial \cdot A=0 . \tag{4.70}
\end{equation*}
$$

Thus, if $m \neq 0$, the constraint $\partial \cdot A=0$ arises classically as a consequence of the equations of motion, not of gauge invariance. In the quantum theory it can indeed be justified to impose $\partial \cdot A=0$ as an operator equation. ${ }^{2}$ The physical Hilbert space now exhibits 3 positive-norm degrees of freedom corresponding to the polarisations

$$
\begin{equation*}
|\vec{k}, \zeta\rangle \text { with } \zeta^{\mu} k_{\mu}=0 \text { and } k^{2}-\mu^{2}=0 . \tag{4.71}
\end{equation*}
$$

Since no residual gauge transformation is available, there is no further decoupling of one degree of freedom. To summarise:

[^21]Massive vector fields have 3 physical degrees of freedom.
Massless vector fields have 2 physical degrees of freedom.

Some comments are in order:

- One might be irritated that the constraint $\partial \cdot A=0$, which in this case rests on the equations of motion, is imposed as a constraint in the quantum theory - quantisation must hold off-shell. For the free theory this is not really a problem: We can take the classical on-shell relation $\partial \cdot A=0$ as a motivation to simply declare the physical Hilbert space to consist of the transverse polarisations only, thereby defining the quantum theory. That this remains correct in the presence of interactions rests again on the Ward identities, which, for suitable interactions, still guarantee that the - now negative norm - states with polarisation $\zeta \sim k$ decouple from the S-matrix. We will discuss this in detail in the context of the Ward identities.
- The procedure for quantising the Proca action breaks down if we set $\mu \rightarrow 0$. A framework where a smooth limit $m \rightarrow 0$ is possible is provided by the Stückelberg Lagrangian ${ }^{3}$

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}+\frac{1}{2} \mu^{2} A_{\mu} A^{\mu}-\frac{\lambda}{2}(\partial A)^{2} . \tag{4.72}
\end{equation*}
$$

### 4.5 Coupling vector fields to matter

We would now like to couple a vector field $A^{\mu}(x)$ to a matter sector, e.g. a Dirac fermion or a scalar field. Let us collectively denote the matter fields by $\phi$ and the pure matter action (excluding any term involving the gauge fields) by $S_{\text {matter }}^{\text {rest }}[\phi]$. Then we would like to construct an action

$$
\begin{equation*}
S=S_{A}^{0}[A]+S_{\mathrm{int}}[A, \phi]+S_{\text {matter }}^{\text {rest }}[\phi] \tag{4.73}
\end{equation*}
$$

that describes the coupling of this matter sector to a vector field $A^{\mu}(x)$ with free action (prior to coupling) $S_{A}^{0}[A]$ and interaction terms $S_{\text {int }}[A, \phi]$. Naively, we would think that we can simply write down all possible Lorentz invariant terms in $S_{\text {int }}[A, \phi]$ involving $A^{\mu}$ and $\phi$ and then organize these as a series in derivatives and powers of fields, as we have done when writing down the most general interaction for a scalar field. This time, however, $S_{\text {int }}[A, \phi]$ is subject to an important constraint: It must be chosen such that the successful decoupling of negative norm states and zero norm states (in the case of a massless vector) in the free vector theory is not spoiled by the interaction. More precisely, if we denote by

$$
\begin{equation*}
\mathcal{M}=\zeta_{\mu}(k) \mathcal{M}^{\mu} \tag{4.74}
\end{equation*}
$$

the scattering amplitude involving an external photon of polarisation vector $\zeta_{\mu}(k)$ and momentum $k$, then consistency of the interaction requires - at the very least - that

$$
\begin{equation*}
k_{\mu} \mathcal{M}^{\mu}=0 \tag{4.75}
\end{equation*}
$$

[^22]This is equivalent to the requirement that external photons of polarisation $\zeta_{\mu}(k)=k_{\mu}$ decouple from the interactions. As we will see this is necessary and also sufficient to ensure that if we start with a physical photon state of positive norm, no negative or zero norm states (in the case of a massless vector theory) are produced via the interaction.
Now, one can (and we will somewhat later in this course) prove very generally the following two important theorems:
1.) The decoupling of unphysical photon states with polarisation $\zeta_{\mu}(k)=k_{\mu}$, i.e. equ. (4.75), is equivalent to the statement that

$$
\begin{equation*}
\frac{\delta S_{\mathrm{int}}[A, \phi]}{\delta A^{\mu}}=-j^{\mu} \tag{4.76}
\end{equation*}
$$

for $j^{\mu}$ the conserved current associated with a global continuous $U(1)$ symmetry of the full action $S$ under which

$$
\begin{equation*}
\phi(x) \rightarrow \phi(x)-e \alpha \delta \phi(x), \quad \alpha \in \mathbb{R} \tag{4.77}
\end{equation*}
$$

infinitesimally. Here we have rescaled the symmetry parameter $\alpha$ by a coupling constant $e$ to comply with later conventions. In particular, $\partial_{\mu} j^{\mu}=0$ on-shell. This is the statement that the vector theory must couple to a conserved current. That coupling to a conserved current is equivalent to (4.75) is ensured by the Ward identities to be discussed in detail later in this course.
2.) If the vector theory is massless, i.e. if $S_{A}^{0}[A]$ is invariant under the gauge symmetry

$$
\begin{equation*}
A_{\mu}(x) \rightarrow A_{\mu}(x)+\partial_{\mu} \alpha(x), \tag{4.78}
\end{equation*}
$$

then (4.76) is equivalent to the statement that the full action $S$ is invariant under the combined gauge transformation

$$
\begin{equation*}
\phi(x) \rightarrow \phi(x)-e \alpha(x) \delta \phi(x) \quad A_{\mu}(x) \rightarrow A_{\mu}(x)+\partial_{\mu} \alpha(x), \tag{4.79}
\end{equation*}
$$

where in particular $\phi(x)$ now transforms under a local symmetry since we have promoted the constant $\alpha$ appearing in (4.77) to a function $\alpha(x)$. This process of promoting the global continuous symmetry (4.77) to a combined gauge symmetry as above is called gauging.

As indicated, for pedagogical reasons we postpone a proof of both these assertions and first exemplify the consistent coupling of a massless vector theory to matter via gauging by discussing interactions with a Dirac fermion and a complex scalar theory.

### 4.5.1 Coupling to Dirac fermions

Let us start with the free Dirac fermion action

$$
\begin{equation*}
S_{\text {matter }}^{\text {rest }}=\int d^{4} x \bar{\psi}\left(i \gamma^{\mu} \partial_{\mu}-m_{0}\right) \psi \tag{4.80}
\end{equation*}
$$

The only available vector Noether current of this action is due to the a priori global $U(1)$ symmetry

$$
\begin{equation*}
\psi(x) \rightarrow e^{-i e \alpha} \psi(x), \quad \bar{\psi}(x) \rightarrow \bar{\psi}(x) e^{i e \alpha} \quad \alpha \in \mathbb{R}, \tag{4.81}
\end{equation*}
$$

where $\alpha \in \mathbb{R}$ is the symmetry parameter and as above we have introduced a dimensionless coupling constant $e$.
With the normalisation conventions of eq. (1.37), the corresponding Noether current is

$$
\begin{equation*}
j^{\mu}=e \bar{\psi} \gamma^{\mu} \psi \text { and } \partial_{\mu} j^{\mu}=0 \tag{4.82}
\end{equation*}
$$

on-shell for $\psi(x)$. As prescribed by the above theorem we proceed by gauging this global $U(1)$ symmetry. This means that we promote the global $U(1)$ symmetry to a local one, i.e. we promote the constant $\alpha \in \mathbb{R}$ to a function $\alpha(x)$. Now the kinetic term is no longer invariant because

$$
\begin{align*}
\bar{\psi}(i \gamma \cdot \partial-m) \psi & \mapsto \bar{\psi} e^{i e \alpha(x)}\left(i \gamma \cdot \partial-m_{0}\right) e^{-i e \alpha(x)} \psi \\
& =\bar{\psi}\left(i \gamma^{\mu}\left(\partial_{\mu}-i e \partial_{\mu} \alpha(x)\right)-m_{0}\right) \psi  \tag{4.83}\\
& =\bar{\psi}\left(i \gamma \cdot \partial-m_{0}\right) \psi+\left[\partial_{\mu} \alpha(x)\right] j^{\mu}(x) .
\end{align*}
$$

However, we observe that the interaction term

$$
\begin{equation*}
\bar{\psi}\left(i \gamma^{\mu} \partial_{\mu}-m_{0}\right) \psi-A_{\mu} j^{\mu}, \quad j^{\mu}=e \bar{\psi} \gamma^{\mu} \psi \tag{4.84}
\end{equation*}
$$

is invariant under the combined gauge transformation

$$
\begin{equation*}
\psi(x) \rightarrow e^{-i e \alpha(x)} \psi(x), \quad A_{\mu} \rightarrow A_{\mu}+\partial_{\mu} \alpha(x) \tag{4.85}
\end{equation*}
$$

Thus the interaction is gauge invariant off-shell and fully consistent. One can rewrite the interaction in terms of the covariant derivative

$$
\begin{equation*}
D_{\mu}:=\partial_{\mu}+i e A_{\mu} . \tag{4.86}
\end{equation*}
$$

Under a combined gauge transformation

$$
\begin{equation*}
\psi(x) \rightarrow e^{-i e \alpha(x)} \psi(x), \quad A_{\mu} \rightarrow A_{\mu}+\partial_{\mu} \alpha(x) \tag{4.87}
\end{equation*}
$$

the covariant derivative transforms as

$$
\begin{equation*}
D_{\mu} \psi(x) \rightarrow e^{-i e \alpha(x)} D_{\mu} \psi(x) \tag{4.88}
\end{equation*}
$$

Thus the Lagrangian

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}+\bar{\psi}\left(i \gamma^{\mu} D_{\mu}-m_{0}\right) \psi=-\frac{1}{4} F^{2}+\bar{\psi}\left(i \gamma \cdot \partial-m_{0}\right) \psi-A_{\mu} j^{\mu}, \quad j^{\mu}=e \bar{\psi} \gamma^{\mu} \psi \tag{4.89}
\end{equation*}
$$

is manifestly gauge invariant off-shell. In QED we set $e=-|e|$ equal to the elementary charge of one electron such that

$$
\begin{align*}
& Q a_{s}^{\dagger}(\vec{p})|0\rangle=-|e| a_{s}^{\dagger}(\vec{p})|0\rangle \text { for an electron, } \\
& Q b_{s}^{\dagger}(\vec{p})|0\rangle=+|e| b_{s}^{\dagger}(\vec{p})|0\rangle \text { for a positron, } \tag{4.90}
\end{align*}
$$

where $Q$ is the Noether charge associated with $j^{\mu}$.

Let us stress that global and gauge symmetries are really on very different footings:

- A gauge symmetry is a redundancy of the description of the system.
- A global symmetry is a true symmetry between different field configurations.

Note furthermore that, by construction, the gauge symmetry (4.87) reduces for $\alpha(x)=\alpha=$ const. to the global symmetry

$$
\begin{equation*}
\psi \rightarrow e^{-i e \alpha} \psi, \quad A \rightarrow A, \tag{4.91}
\end{equation*}
$$

from which, in turn, charge conservation follows. In particular the appearance of a combined $U(1)$ gauge symmetry of the system matter plus gauge fields requires an underlying global $U(1)$ symmetry of the matter system.

### 4.5.2 Coupling to scalars

Conserved vector currents are available only for complex scalars. Put differently, real scalars are uncharged. Let us therefore consider a complex scalar theory, whose action prior to coupling to the Maxwell field reads

$$
\begin{equation*}
S_{\mathrm{matter}}^{\mathrm{rest}}=\partial_{\mu} \phi^{\dagger}(x) \partial^{\mu} \phi(x)-m^{2} \phi^{\dagger}(x) \phi(x) \tag{4.92}
\end{equation*}
$$

The Noether current of the free theory associated with the global $U(1)$ symmetry $\phi \rightarrow e^{-i e \alpha} \phi$ with $\alpha \in \mathbb{R}$ is $j_{\text {free }}^{\mu}=i e\left(\phi^{\dagger} \partial^{\mu} \phi-\left(\partial^{\mu} \phi^{\dagger}\right) \phi\right)$.
The naive guess for the coupling to the gauge sector would therefore be

$$
\begin{equation*}
\mathcal{L}_{\text {int }}^{\text {naive }}=-A_{\mu} j_{\text {free }}^{\mu} . \tag{4.93}
\end{equation*}
$$

However - unlike in the fermionic case - this coupling does not exhibit off-shell gauge invariance under

$$
\begin{equation*}
\phi \rightarrow e^{-i e \alpha(x)} \phi, \quad A_{\mu} \rightarrow A_{\mu}+\partial_{\mu} \alpha(x) . \tag{4.94}
\end{equation*}
$$

Let us instead follow the general route of replacing the usual derivative $\partial_{\mu} \phi$ by the covariant derivative

$$
\begin{equation*}
D_{\mu} \phi(x)=\left(\partial_{\mu}+i e A_{\mu}\right) \phi . \tag{4.95}
\end{equation*}
$$

and consider the action

$$
\begin{equation*}
S=\int d^{4} x\left(-\frac{1}{4} F^{2}+\left(D_{\mu} \phi\right)^{\dagger} D^{\mu} \phi-m^{2} \phi^{\dagger} \phi\right) . \tag{4.96}
\end{equation*}
$$

The interaction which follows by expanding $D^{\mu}=\partial^{\mu}+i e A^{\mu}$ is

$$
\begin{equation*}
S_{\mathrm{int}}[A, \phi]=-\int d^{4} x\left(i e\left(\phi^{\dagger} \partial^{\mu} \phi-\left(\partial^{\mu} \phi^{\dagger}\right) \phi\right) A_{\mu}-e^{2} A_{\mu} A^{\mu} \phi^{\dagger} \phi\right) . \tag{4.97}
\end{equation*}
$$

The last term quadratic in $A^{\mu}$ is required for gauge invariance and was missed in the naive guess (4.93).

To see what had gone wrong in the naive guess (4.93) note that for constant gauge parameter $\alpha \in \mathbb{R}$ the combined gauge transformation (4.94) reduces to a $U(1)$ global symmetry of the full action (4.96). One can check that the Noether current associated with this global $U(1)$ symmetry of (4.96) is just

$$
\begin{equation*}
j^{\mu}=-2 e^{2} A^{\mu} \phi^{\dagger} \phi+i e\left(\phi^{\dagger} \partial^{\mu} \phi-\left(\partial^{\mu} \phi^{\dagger}\right) \phi\right)=i e\left(\phi^{\dagger} D^{\mu} \phi-\left(D^{\mu} \phi\right)^{\dagger} \phi\right) \tag{4.98}
\end{equation*}
$$

and this does agree with the formula

$$
\begin{equation*}
\frac{\delta S_{\mathrm{int} .}[A, \phi]}{\delta A_{\mu}}=-j^{\mu} \tag{4.99}
\end{equation*}
$$

The point is that in general the Noether current may itself depend on $A_{\mu}$ once the coupling to the gauge field is taken into account because $S_{\text {int. }}[A, \phi]$ may depend not just linearly on $A$. Therefore merely writing $-A_{\mu} j_{\text {free }}^{\mu}$ with $j_{\text {free }}^{\mu}$ the conserved current associated with $S_{\text {matter }}^{\text {rest }}$ only is in general not the right thing to do.
We conclude that generally the correct way to define gauge invariant interactions is by replacing $\partial_{\mu} \rightarrow D_{\mu}$. This is called minimal coupling. In section 5.2 we will show that gauge invariance of the combined matter and gauge sector is necessary and sufficient to define a consistent theory with in particular consistent interactions.

### 4.6 Feynman rules for QED

We are finally in a position to study the interactions of Quantum Electrodynamics (QED), whose Lagrangian is given by

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}-\frac{\lambda}{2}(\partial \cdot A)^{2}+\bar{\psi}\left(i \gamma \cdot \partial-m_{0}\right) \psi-e A_{\mu} \bar{\psi} \gamma^{\mu} \psi \tag{4.100}
\end{equation*}
$$

This theory describes the coupling of the Maxwell $U(1)$ gauge potential to electro-magnetically charged spin $\frac{1}{2}$ particles of free mass $m_{0}$.
The Feynman rules for the $U(1)$ gauge field are simple to state:

- The Feynman propagator for the gauge field in Feynman gauge $(\lambda=1)$ can easily be computed as

$$
\begin{equation*}
\langle 0| T A^{\mu}(x) A^{v}(y)|0\rangle=-\left.\eta^{\mu v} D_{F}^{(0)}(x-y)\right|_{m_{0}^{2}=0}=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{-i \eta^{\mu v}}{p^{2}+i \epsilon} e^{-i p \cdot(x-y)} \tag{4.101}
\end{equation*}
$$

by plugging in the mode expansion for the quantised spin 1 field in the Heisenberg picture and proceeding as in the scalar field case.

- For completeness we note that for arbitrary $\lambda$ the propagator is

$$
\begin{equation*}
\langle 0| T A^{\mu}(x) A^{v}(y)|0\rangle=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{-i\left(\eta^{\mu \nu}+\frac{p^{\mu} p^{v}}{p^{2}}\left(\frac{1}{\lambda}-1\right)\right)}{p^{2}+i \epsilon} e^{-i p \cdot(x-y)} \tag{4.102}
\end{equation*}
$$

This is proven most easily in path-integral quantisation as will be discussed in detail in the course QFT II.

- Graphically we represent the propagator of a gauge field as follows:

$$
{ }^{y} \sim \sim \sim \sim^{x} \equiv-\eta_{\mu \nu} D_{F}^{(0)}(x-y) .
$$

- To determine the Feynman rules we must go through the LSZ analysis for gauge fields. By arguments similar to the ones that lead to the appearance of the spinor polarisation in the spin $1 / 2$ case one finds that external photon states $|\vec{p}, \lambda\rangle$ come with polarisation factors

$$
\left.\begin{array}{cc}
\epsilon_{\mu}(\vec{p}, \lambda) & \text { for ingoing }  \tag{4.103}\\
\epsilon_{\mu}^{*}(\vec{p}, \lambda) & \text { for outgoing }
\end{array}\right\}|\vec{p}, \lambda\rangle
$$

Here we allowed for complex polarisation vectors. ${ }^{4}$
This leads to the following Feynman rules of QED:
To compute the scattering amplitude $i \mathcal{M}_{f i}$ defined in equ. (2.165) of a given process we draw all relevant fully connected, amputated Feynman diagrams to the given order in the coupling constant $e$ and read off $i \mathcal{M}_{f i}$ as follows:

- Each interaction vertex has the form

(with the arrows denoting fermion number flow) and carries a factor $-i e \gamma^{\mu}$,
- each internal photon line $\mu \sim \sim \sim n \sim \sim v$ carries a factor $-\frac{i \eta^{\mu \nu}}{p^{2}+i \epsilon}$,
- each internal fermion line $\longrightarrow$ with the arrow denoting fermion (as opposed to antifermion) number flow carries a factor $\frac{i\left(\gamma \cdot p+m_{0}\right)}{p^{2}-m_{0}^{2}+i \epsilon}$,
- momentum conservation is imposed at each vertex,
- we integrate over each (undetermined) internal momentum with the usual measure $\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}}$,
- each ingoing photon of polarisation $\lambda$ carries a factor $\epsilon^{\mu}(\vec{p}, \lambda) Z_{A}^{\frac{1}{2}}$,
each outgoing photon of polarisation $\lambda$ carries a factor $\epsilon^{\mu *}(\vec{p}, \lambda) Z_{A}^{\frac{1}{2}}$,
- each ingoing fermion of spin $s$ carries a factor $u_{s}(\vec{p}) Z_{e}^{\frac{1}{2}}$,
each ingoing anti-fermion of spin $s$ carries a factor $\bar{v}_{-s}(\vec{p}) Z_{e}^{\frac{1}{2}}$,

[^23]- each outgoing fermion of spin $s$ carries a factor $\bar{u}_{s}(\vec{p}) Z_{e}^{\frac{1}{2}}$,
each outgoing anti-fermion of spin $s$ carries a factor $v_{-s}(\vec{p}) Z_{e}^{\frac{1}{2}}$.
- The overall sign of a given diagram is easiest determined directly in the interaction picture. If 2 diagrams are related by the exchange of $n$ fermion lines, then the relative sign is $(-1)^{n}$. If we are only interested in $\left|i \mathcal{M}_{f i}\right|^{2}$ this is often enough to determine the cross-section.

We will now give some important examples of leading order QED processes:

## Electron scattering

Electron-scattering corresponds to the process

$$
\begin{equation*}
e^{-} e^{-} \rightarrow e^{-} e^{-} \tag{4.104}
\end{equation*}
$$

To leading order in the coupling $e$ this process receives contributions from the two following fully connected, amputated diagrams:


Figure 4.1: Leading order Feynman diagrams for electron scattering.

According to the Feynman rules

$$
\begin{align*}
i \mathcal{M}_{f i}= & (i e)^{2}\left[\bar{u}_{r^{\prime}}\left(q^{\prime}\right) \gamma^{\mu} u_{r}(q) \frac{-i \eta_{\mu \nu}}{t^{2}+i \epsilon} \bar{u}_{s^{\prime}}\left(p^{\prime}\right) \gamma^{v} u_{s}(p)\right.  \tag{4.105}\\
& \left.-\bar{u}_{s^{\prime}}\left(p^{\prime}\right) \gamma^{\mu} u_{r}(q) \frac{-i \eta_{\mu \nu}}{t^{\prime 2}+i \epsilon} \bar{u}_{r^{\prime}}\left(q^{\prime}\right) \gamma^{v} u_{s}(p)\right]
\end{align*}
$$

where

$$
\begin{equation*}
q-t-q^{\prime}=0 \Rightarrow t=q-q^{\prime} \text { and } t+p-p^{\prime}=0 \Rightarrow p^{\prime}=p+q-q^{\prime} \tag{4.106}
\end{equation*}
$$

in the first diagram and

$$
\begin{equation*}
t^{\prime}=q-p^{\prime} \text { and } q^{\prime}=q-p^{\prime}+p \tag{4.107}
\end{equation*}
$$

in the second diagram. Note the relative minus sign between both diagrams due to the crossing fermion line!

## Electron-positron annihilation

Electron-positron annihilation is the process

$$
\begin{equation*}
e^{+} e^{-} \rightarrow 2 \gamma \tag{4.108}
\end{equation*}
$$

The two leading-order diagrams are


Figure 4.2: Electron-positron annihilation leading-order diagrams.

Then the Feynman rules yield

$$
\begin{align*}
i \mathcal{M}_{f i}= & (i e)^{2}\left[\bar{v}_{-r}(q) \gamma^{\nu} \frac{i\left(\gamma t-m_{0}\right)}{t^{2}-m_{0}^{2}+i \epsilon} \gamma^{\mu} u_{s}(p) \epsilon_{\nu}^{*}\left(q^{\prime}, \lambda^{\prime}\right) \epsilon_{\mu}^{*}\left(p^{\prime}, \lambda\right)\right.  \tag{4.109}\\
& \left.+\bar{v}_{-r}(q) \gamma^{\nu} \frac{i\left(\gamma t^{\prime}+m_{0}\right)}{t^{\prime 2}-m_{0}^{2}+i \epsilon} \gamma^{\mu} u_{s}(p) \epsilon_{\nu}^{*}\left(p^{\prime}, \lambda\right) \epsilon_{\mu}^{*}\left(q^{\prime}, \lambda^{\prime}\right)\right],
\end{align*}
$$

where momentum conservation implies

$$
\begin{equation*}
t=p-p^{\prime}, \quad t^{\prime}=p-q^{\prime}, \quad q^{\prime}=p-p^{\prime}+q, \quad p^{\prime}=p-q^{\prime}+q . \tag{4.110}
\end{equation*}
$$

## $e^{+} e^{-}$- scattering (Bhabha scattering)

The process

$$
\begin{equation*}
e^{+} e^{-} \rightarrow e^{+} e^{-} \tag{4.111}
\end{equation*}
$$

is described by the leading order diagram 4.3.

## Compton scattering

Compton scattering corresponds to

$$
\begin{equation*}
e^{-}+\gamma \rightarrow e^{-}+\gamma \tag{4.112}
\end{equation*}
$$

with leading order diagram 4.4.


Figure 4.3: Bhabha scattering at leading order.


Figure 4.4: Compton scattering at leading order.

## Non-linearities

As an interesting new quantum effect, loop diagrams induce an interaction between photons. For example, at one loop two photons scatter due to a process of the form


Figure 4.5: Photon scattering is not allowed in classical Electrodynamics, but it is in QED.

Such effects are absent in the classical theory, where light waves do not interact with each other due to the linear structure of the classical theory. It is a fascinating phenomenon that such QED effects break the linearity of classical optics. In high intensity laser beams, where quantum effects are quantitatively relevant, such non-linear optics phenomena can indeed be observed.

### 4.7 Recovering Coulomb's potential

As we have seen QED interactions are mediated by the exchange of gauge bosons. It is an interesting question to determine the physical potential $V(\vec{x})$ induced by the exchange of such gauge bosons. This will teach us how the concept of classical forces emerges from the QFT framework of scattering.

To derive $V(\vec{x})$ the key idea is to consider the non-relativistic limit of the electron scattering process

$$
\begin{equation*}
e^{-}+e^{-} \rightarrow e^{-}+e^{-} \tag{4.113}
\end{equation*}
$$

with distinguishable out-states corresponding to the Feynman diagram 4.6.


Figure 4.6: Electron scattering with distinguishable out-states.

This is then compared to the non-relativistic scattering process $|\vec{p}\rangle \rightarrow|\vec{p}\rangle$ of two momentum eigenstates off a potential $V(\vec{x})$.

- In non-relativistic quantum mechanics, the scattering amplitude $\mathcal{A}\left(|\vec{p}\rangle \rightarrow\left|\vec{p}^{\prime}\right\rangle\right)$ for scattering of an incoming momentum eigenstate $|\vec{p}\rangle$ to an outgoing momentum eigenstate $\left|\vec{p}^{\prime}\right\rangle$ in the presence of a scattering potential $V$ is computed in the interaction picture as

$$
\mathcal{A}\left(|\vec{p}\rangle \rightarrow\left|\vec{p}^{\prime}\right\rangle\right)-1=\left\langle\vec{p}^{\prime}\right| e^{-i \int_{-\infty}^{\infty} \mathrm{d} t^{\prime} V_{I}\left(t^{\prime}\right)}|\vec{p}\rangle-1 \cong-i\langle\vec{p}| \int_{-\infty}^{\infty} \mathrm{d} t^{\prime} V_{I}\left(t^{\prime}\right)|\vec{p}\rangle+\ldots,
$$

where we take the potential to be constant in time. As derived in standard textbooks on Quantum Mechanics this equals to first order

$$
\begin{equation*}
\mathcal{A}\left(|\vec{p}\rangle \rightarrow\left|\vec{p}^{\prime}\right\rangle\right)-1=(2 \pi) \delta\left(E_{p}-E_{p^{\prime}}\right)(-i)\left\langle\vec{p}^{\prime}\right| V|\vec{p}\rangle . \tag{4.114}
\end{equation*}
$$

This result goes by the name Born approximation and reads more explicitly in position space

$$
\begin{equation*}
\mathcal{A}\left(|\vec{p}\rangle \rightarrow\left|\vec{p}^{\prime}\right\rangle\right)-1=(2 \pi) \delta\left(E_{p}-E_{p^{\prime}}\right)(-i) \int \mathrm{d}^{3} r V(\vec{r}) e^{-i\left(\vec{p}-\vec{p}^{\prime}\right) \cdot \vec{r}} \tag{4.115}
\end{equation*}
$$

- This is to be compared with the non-relativistic limit of

$$
\begin{equation*}
\left.\int \frac{\mathrm{d}^{3} k}{(2 \pi)^{3}} e^{i \vec{k} \cdot \vec{r}_{0}}\left\langle p^{\prime}, k^{\prime}\right| S|p, k\rangle\right|_{\text {conn. }} \tag{4.116}
\end{equation*}
$$

where $\left.\left\langle p^{\prime}, k^{\prime}\right| S|p, k\rangle\right|_{\text {conn. }}$ is the connected part of the S-matrix element of the scattering process (4.113). Here the idea is to identify the electron with momentum $p$ and $p^{\prime}$ with the scattering states in the Quantum Mechanics picture and to view the other electron as a fixed target at $\vec{r}_{0}$ off which $|\vec{p}\rangle$ scatters. The localisation of the fixed target at $\vec{r}_{0}$ in space is achieved by integrating over all Fourier modes, hence explaining the factor $\int \frac{d^{3} k}{(2 \pi)^{3}} e^{i \vec{k} \cdot \vec{r}_{0}}$. Without loss of generality we will set $\vec{r}_{0} \equiv 0$ in the sequel. Recalling the connection between the S-matrix element and the scattering amplitude $i \mathcal{M}$, for which our Feynman rules are formulated, we find

$$
\begin{align*}
\left.\left\langle p^{\prime}, k^{\prime}\right| S|p, k\rangle\right|_{\text {conn. }}= & (2 \pi)^{4} \delta^{(4)}\left(p_{f}-p_{i}\right)(-i e)^{2}  \tag{4.117}\\
& \times\left[\bar{u}_{r^{\prime}}\left(p^{\prime}\right) \gamma^{\mu} u_{r}(p) \frac{-i \eta_{\mu \nu}}{\left(p-p^{\prime}\right)^{2}+i \epsilon} \bar{\epsilon}_{s^{\prime}}\left(k^{\prime}\right) \gamma^{v} u_{s}(k)\right] .
\end{align*}
$$

- In the the non-relativistic limit we have $m_{0} \gg|\vec{p}|$ and $m_{0} \gg|\vec{k}|$ and thus approximate $p \cong$ $\left(m_{0}, \vec{p}\right)+O\left(\vec{p}^{2}\right)$. Then

$$
\begin{equation*}
u_{r}(\vec{p})=\binom{\sqrt{\sigma^{\mu} p_{\mu}} \xi_{r}}{\sqrt{\sigma^{\mu} p_{\mu}} \xi_{r}} \rightarrow\binom{\sqrt{\sigma^{0} m_{0}} \xi_{r}}{\sqrt{\sigma^{0} m_{0}} \xi_{r}}=\sqrt{m_{0}}\binom{\xi_{r}}{\xi_{r}} \tag{4.118}
\end{equation*}
$$

and

$$
\bar{u}_{r^{\prime}}\left(\vec{p}^{\prime}\right) \gamma^{\mu} u_{r}(\vec{p}) \rightarrow\left\{\begin{array}{cc}
2 m_{0} \delta_{r r^{\prime}} & \text { if } \mu=0,  \tag{4.119}\\
0 & \text { otherwise } .
\end{array}\right.
$$

Furthermore $\left(p-p^{\prime}\right)^{2}=-\left|\vec{p}-\vec{p}^{\prime}\right|^{2}+O\left(|\vec{p}|^{4}\right)$ and therefore in the non-relativistic limit

$$
\begin{equation*}
\left.\left\langle p^{\prime}, k^{\prime}\right| S|p, k\rangle\right|_{\text {conn. }} \simeq-i \frac{e^{2}}{\left|\vec{p}-\vec{p}^{\prime}\right|^{2}-i \epsilon}(2 m)^{2} \delta\left(E_{f}-E_{i}\right)(2 \pi)^{4} \delta^{(3)}\left(\vec{p}_{f}-\vec{p}_{i}\right) . \tag{4.120}
\end{equation*}
$$

The factor $(2 m)^{2}$ is due to the different normalisation of momentum eigenstates in QFT and in QM and must therefore be neglected in comparing the respective expressions for the amplitude.

Putting all pieces together we can now identify

$$
\begin{equation*}
\int \mathrm{d}^{3} r V(\vec{r}) e^{-i\left(\vec{p}-\vec{p}^{\prime}\right) \cdot \vec{r}}=\frac{e^{2}}{\left|\vec{p}-\vec{p}^{\prime}\right|^{2}-i \epsilon} . \tag{4.121}
\end{equation*}
$$

Thus the Coulomb potential is

$$
\begin{align*}
V(\vec{r}) & =\int \frac{\mathrm{d}^{3} q}{(2 \pi)^{3}} \frac{e^{2}}{\vec{q}^{2}-i \epsilon} e^{i} e^{\equiv q r \cos \cdot \vec{q} \cdot \vec{r}} \tag{4.122}
\end{align*}=\frac{e^{2}}{4 \pi^{2}} \int_{0}^{\infty} \mathrm{d} q \frac{q^{2}}{q^{2}-i \epsilon} \frac{e^{i q r}-e^{-i q r}}{i q r}
$$

Note that $\sqrt{i \epsilon}=e^{i \pi / 4} \sqrt{\epsilon}$ lies in the upper complex half-plane. We interpret the integral as a complex contour integral in the complex upper half-plane because we can close the contour in the upper halfplane as the contribution above the real axis vanishes in the limit $|q| \rightarrow \infty$. We then pick up the residue at $q=\sqrt{+i \epsilon}$ and afterwards take the limit $\epsilon \rightarrow 0$. This leads to

$$
\begin{equation*}
V(\vec{r})=\left.\frac{e^{2}}{4 \pi^{2} i r} 2 \pi i \frac{\sqrt{i \epsilon}}{2 \sqrt{i \epsilon}} e^{i \sqrt{i \epsilon r}}\right|_{\epsilon=0}=\frac{e^{2}}{4 \pi r} . \tag{4.123}
\end{equation*}
$$

Since $V$ is positive the potential for scattering of two electrons is repulsive as expected. If we replace one $e^{-}$by $e^{+}$and consider instead the process

$$
\begin{equation*}
e^{-}+e^{+} \rightarrow e^{-}+e^{+} \tag{4.124}
\end{equation*}
$$

this amounts to replacing

$$
\bar{u}_{r^{\prime}}\left(p^{\prime}\right) \gamma^{\mu} u_{r}(p) \text { by } \bar{v}_{r^{\prime}}\left(p^{\prime}\right) \gamma^{\mu} v_{r}(p) \rightarrow\left\{\begin{array}{cl}
2 m_{0} \delta_{r r^{\prime}} & \text { if } \mu=0  \tag{4.125}\\
0 & \text { otherwise }
\end{array}\right.
$$

However, a careful analysis of the scattering amplitude in the interaction picture shows that one picks up one relative minus sign in the amplitude. Indeed, establishing the absolute sign of the amplitude will be the task of Assignment 11. Thus we confirm the expected result that the potential mediated by exchange of a spin-1 particle yields an attractive potential for oppositely charged particles and a repulsive potential for identically charged particles.

### 4.7.1 Massless and massive vector fields

It is instructive to compare the interaction of massless and massive vector fields:

- The Coulomb potential due to exchange of a massless spin-1 particle leads to a long-range force as the potential dies off only like $1 / r$. If on the other hand the exchanged spin-1 particle is massive, the interaction is short-ranged. To see this suppose the photon has mass $\mu$. All that changes is a mass term in the phton propagator. Thus

$$
\begin{equation*}
V(\vec{r})=\int \frac{\mathrm{d}^{3} q}{(2 \pi)^{3}} \frac{e^{2}}{\vec{q}^{2}+\mu^{2}-i \epsilon} e^{i \vec{q} \cdot \vec{r}}=\frac{e^{2}}{4 \pi r} e^{-\mu r} . \tag{4.126}
\end{equation*}
$$

## Therefore massive vector bosons lead to short-range forces of range $\sim \frac{1}{\mu}$.

While the vector boson of QED is massless ${ }^{5}$, the weak nuclear forces in the Standard Model are mediated by massive vector bosons (the three massive vector bosons $W^{+}, W^{-}$and $Z$ of spontaneously broken $S U(2)$ gauge symmetry of mass of order 100 GeV ). So this effect is

[^24]indeed highly relevant in particle physics and explains why in everyday physics at distances bigger than $(100 \mathrm{GeV})^{-1}$ only the electromagnetic force can be experienced. ${ }^{6}$

- The concept of forces mediated by exchange bosons is not restricted to spin-1 theories.

If we replace the vector boson of QED by a scalar boson $\phi$ of mass $\mu$ we arrive at Yukawa theory defined by

$$
\begin{equation*}
\mathcal{L}=\frac{1}{2} \partial \phi^{2}-\frac{1}{2} \mu^{2} \phi^{2}+\bar{\psi}\left(i \gamma^{\mu} \partial_{\mu}-m_{0}\right) \psi-\underbrace{e \phi \bar{\psi} \psi}_{\text {Yukawa int. }} . \tag{4.127}
\end{equation*}
$$

The Feynman rules are very similar to those in QED except for two important changes: First, the interaction vertex carries no $\gamma^{\mu}$ factor and second the scalar boson propagator must be modified from

$$
\begin{equation*}
\frac{-i \eta_{\mu \nu}}{p^{2}-\mu^{2}+i \epsilon} \rightarrow \frac{i}{p^{2}-\mu^{2}+i \epsilon} . \tag{4.128}
\end{equation*}
$$

This amounts to a crucial sign change in the non-relativistic limit because

$$
\begin{equation*}
\frac{-i \eta_{00}}{-|\vec{p}|^{2}-\mu^{2}+i \epsilon} \rightarrow \frac{i}{-|\vec{p}|^{2}-\mu^{2}+i \epsilon} . \tag{4.129}
\end{equation*}
$$

Combining these two changes yields the universally attractive Yukawa potential

$$
\begin{equation*}
V(r)=-\frac{e^{2}}{4 \pi r} e^{-\mu r} . \tag{4.130}
\end{equation*}
$$

- Perturbative gravity can be understood at the level of field theory as a theory of massless spin-2 particles, the gravitons, whose exchange likewise yields a universally attractive force. Since the gravitational potential is of the form $1 / r$ (at least for conventional Einstein gravity), gravitons must be massless. ${ }^{7}$

To summarise the potential induced by the exchange of bosons of different spin acts on fermions ( $f$ ) and anti-fermions $(\bar{f})$ as follows:

|  | $f f$ | $f \bar{f}$ | $\bar{f} \bar{f}$ |
| :--- | :---: | :---: | :---: |
| Spin 0 (Yukawa theory) | attractive | attractive | attractive |
| Spin 1 (Gauge theory) | repulsive | attractive | repulsive |
| Spin 2 (Gravity) | attractive | attractive | attractive |

[^25]
## Chapter 5

## Quantum Electrodynamics

### 5.1 QED process at tree-level

As an example for a typical tree-level QED process we compute the differential cross-section $\frac{\mathrm{d} \sigma}{\mathrm{d} \Omega}$ for the scattering $e^{+} e^{-} \rightarrow \mu^{+} \mu^{-}$. We describe both the electron/positron and the muon/anti-muon by a Dirac spinor field of respective free mass $m_{e}$ and $m_{\mu}$ which we couple to the $U(1)$ gauge field.


Figure 5.1: $e^{+} e^{-} \rightarrow \mu^{+} \mu^{-}$reaction.

In this protypical example and in many similar processes one proceeds as follows:

### 5.1.1 Feynman rules for in/out-states of definite polarisation

Apply the rules from the previous chapter and obtain:

$$
\begin{equation*}
i \mathcal{M}\left(p, s ; p^{\prime},-s^{\prime} \rightarrow k, r ; k^{\prime},-r^{\prime}\right)=(i e)^{2} \bar{v}_{s^{\prime}}\left(p^{\prime}\right) \gamma^{\mu} u_{s}(p) \frac{-i \eta_{\mu v}}{q^{2}+i \epsilon} \bar{u}_{r}(k) \gamma^{v} v_{r^{\prime}}\left(k^{\prime}\right), \tag{5.1}
\end{equation*}
$$

where $q=p+p^{\prime}=k+k^{\prime}$. Thus

$$
\begin{equation*}
i \mathcal{M}=\frac{i e^{2}}{q^{2}} \bar{v}_{s^{\prime}}\left(p^{\prime}\right) \gamma^{\mu} u_{s}(p) \bar{u}_{r}(k) \gamma_{\mu} v_{r^{\prime}}\left(k^{\prime}\right) \tag{5.2}
\end{equation*}
$$

### 5.1.2 Sum over all spin and polarisation states

Often we do not keep track of the polarisation states of in and out states but are interested in the unpolarised amplitude-square

$$
\begin{equation*}
\frac{1}{2} \sum_{s} \frac{1}{2} \sum_{s^{\prime}} \sum_{r} \sum_{r^{\prime}}\left|i \mathcal{M}\left(p, s ; p^{\prime},-s^{\prime} \rightarrow k, r ; k^{\prime},-r^{\prime}\right)\right|^{2}, \tag{5.3}
\end{equation*}
$$

where $\frac{1}{2} \sum_{s} \frac{1}{2} \sum_{s^{\prime}}$ averages over the initial state polarisation as is appropriate if the incoming beam is not prepared in a polarisation eigenstate. The sum $\sum_{r} \sum_{r^{\prime}}$ over the final state polarisations is required in addition if the detector is blind to polarisation. With

$$
\begin{equation*}
\left(\bar{v}_{s^{\prime}} \gamma^{\mu} u_{s}\right)^{*}=u_{s}^{\dagger} \underbrace{\gamma^{\dagger \dagger} \gamma^{0 \dagger}}_{=\gamma^{0} \gamma^{\mu}} v_{s^{\prime}}=\bar{u}_{s} \gamma^{\mu} v_{s^{\prime}} \tag{5.4}
\end{equation*}
$$

one finds

$$
\begin{equation*}
\frac{1}{4} \sum_{\text {Spins }}|\mathcal{M}|^{2}=\frac{1}{4} \frac{e^{2}}{q^{4}}\left[\sum_{s, s^{\prime}, r, r^{\prime}}\left(\bar{u}_{s} \gamma^{\mu} v_{s^{\prime}}\right)\left(\bar{v}_{s^{\prime}} \gamma^{v} u_{s}\right)\left(\bar{v}_{r^{\prime}} \gamma_{\mu} u_{r}\right)\left(\bar{u}_{r} \gamma_{\nu} v_{r^{\prime}}\right)\right] . \tag{5.5}
\end{equation*}
$$

Now we want to make use of the completeness relations

$$
\begin{align*}
& \sum_{s} u_{s}(\vec{p}) \bar{u}_{s}(\vec{p})=\gamma \cdot p+m_{e} \\
& \sum_{s} v_{s}(\vec{p}) \bar{v}_{s}(\vec{p})=\gamma \cdot p-m_{\mu} \tag{5.6}
\end{align*}
$$

To do so it is first useful to make the spinor indices explicit ${ }^{1}$, e.g.

$$
\begin{equation*}
\bar{u}_{s} \gamma^{\mu} v_{s^{\prime}}=\sum_{A, B}\left(\bar{u}_{s}\right)_{A} \gamma^{\mu A}{ }_{B} v_{s^{\prime}}{ }^{B}, \tag{5.7}
\end{equation*}
$$

which shows that the expression can be reordered as

$$
\begin{align*}
\frac{1}{4} \sum_{\text {Spins }}|\mathcal{M}|^{2}= & \frac{1}{4} \frac{e^{2}}{q^{4}}\left[\sum_{s, s^{\prime}, r, r^{\prime}} u_{s}(p)^{D} \bar{u}_{s}(p)_{A} \gamma^{\mu A}{ }_{B} v_{s^{\prime}}\left(p^{\prime}\right)^{B} \bar{v}_{s^{\prime}}\left(p^{\prime}\right)_{C} \gamma^{\nu C}{ }_{D}\right. \\
& \left.\times v_{r^{\prime}}\left(k^{\prime}\right)^{H} \bar{v}_{r^{\prime}}\left(k^{\prime}\right)_{E} \gamma_{\mu}{ }^{E}{ }_{F} u_{r}(k)^{F} \bar{u}_{r}(k)_{G} \gamma_{v}{ }_{H}{ }_{H}\right]  \tag{5.8}\\
= & \frac{1}{4} \frac{e^{2}}{q^{4}}\left[\operatorname{tr}\left(\left(\gamma \cdot p-m_{e}\right) \gamma^{\mu}\left(\gamma \cdot p+m_{e}\right) \gamma^{v}\right)\right. \\
& \left.\times \operatorname{tr}\left(\left(\gamma \cdot k^{\prime}+m_{\mu}\right) \gamma_{\mu}\left(\gamma \cdot k-m_{\mu}\right) \gamma_{\nu}\right)\right] .
\end{align*}
$$

In order to proceed further we need to evaluate the traces over the $\gamma$ matrices.

[^26]
### 5.1.3 Trace identities

The following identities are very important to proceed further. We will prove them in general fashion on exercise sheet 10 . The only thing to do is: insert a $\mathbb{1}$ cleverly, e.g. $\gamma^{5} \gamma^{5}$ and use the properties of the Clifford algebra as well as the cyclicity of the trace.

- $\operatorname{tr} \gamma^{\mu}=0$, because:

$$
\begin{equation*}
\operatorname{tr} \gamma^{\mu}=\operatorname{tr} \underbrace{\gamma^{5} \gamma^{5}}_{\equiv \mathbb{1}} \gamma^{\mu}=-\operatorname{tr} \gamma^{5} \gamma^{\mu} \gamma^{5}=-\operatorname{tr} \gamma^{\mu} \gamma^{5} \gamma^{5}=-\operatorname{tr} \gamma^{\mu} . \tag{5.9}
\end{equation*}
$$

- The trace of an odd number of $\gamma$-matrices vanishes, since

$$
\begin{equation*}
\operatorname{tr} \gamma^{\mu_{1}} \ldots \gamma^{\mu_{n}}=\operatorname{tr} \gamma^{5} \gamma^{5} \gamma^{\mu_{1}} \ldots \gamma^{\mu_{n}}=(-1)^{n} \operatorname{tr} \gamma_{1}^{\mu} \ldots \gamma^{\mu_{n}} . \tag{5.10}
\end{equation*}
$$

- $\operatorname{tr} \gamma^{\mu} \gamma^{\nu}=4 \eta^{\mu \nu}$, because

$$
\begin{equation*}
\operatorname{tr} \gamma^{\mu} \gamma^{\nu}=\operatorname{tr}\left(2 \eta^{\mu \nu}-\gamma^{\nu} \gamma^{\mu}\right)=2 \eta^{\mu \nu} \operatorname{tr} \mathbb{1}-\underbrace{\operatorname{tr}\left(\gamma^{\nu} \gamma^{\mu}\right)}_{=\operatorname{tr} \gamma^{\mu} \gamma^{\nu}} . \tag{5.11}
\end{equation*}
$$

- Similarly one finds

$$
\begin{equation*}
\operatorname{tr} \gamma^{\mu} \gamma^{\nu} \gamma^{\rho} \gamma^{\sigma}=4\left[\eta^{\mu \nu} \eta^{\rho \sigma}-\eta^{\mu \rho} \eta^{\nu \sigma}+\eta^{\mu \sigma} \eta^{\nu \rho}\right] . \tag{5.12}
\end{equation*}
$$

- $\operatorname{tr} \gamma^{5}=0$ since

$$
\begin{equation*}
\operatorname{tr} \gamma^{5}=\operatorname{tr} \gamma^{0} \gamma^{0} \gamma^{5}=-\operatorname{tr} \gamma^{0} \gamma^{5} \gamma^{0}=-\operatorname{tr} \gamma^{5} \tag{5.13}
\end{equation*}
$$

- Likewise one can prove

$$
\begin{equation*}
\operatorname{tr} \gamma^{5} \gamma^{\mu}=\operatorname{tr} \gamma^{5} \gamma^{\mu} \gamma^{\nu}=\operatorname{tr} \gamma^{5} \gamma^{\mu} \gamma^{v} \gamma^{\sigma}=0 . \tag{5.14}
\end{equation*}
$$

- $\operatorname{tr} \gamma^{5} \gamma^{\alpha} \gamma^{\beta} \gamma^{\gamma} \gamma^{\delta}=-4 i \epsilon^{\alpha \beta \gamma \delta}$, because the result must be antisymmetric in all indices and in particular

$$
\begin{equation*}
\operatorname{tr} \gamma^{5} \gamma^{0} \gamma^{1} \gamma^{2} \gamma^{3}=-i \operatorname{tr} \gamma^{5} \gamma^{5}=-4 i \tag{5.15}
\end{equation*}
$$

- Further useful identities are

$$
\begin{align*}
\gamma^{\mu} \gamma^{\nu} \gamma_{\mu} & =-2 \gamma^{v} . \\
\gamma^{\mu} \gamma^{v} \gamma^{\rho} \gamma_{\mu} & =4 \eta^{v \sigma},  \tag{5.16}\\
\gamma^{\mu} \gamma^{\nu} \gamma^{\rho} \gamma^{\sigma} \gamma_{\mu} & =-2 \gamma^{\sigma} \gamma^{\rho} \gamma^{\nu} .
\end{align*}
$$

Applied to the present case these identities yield

$$
\begin{align*}
& \operatorname{tr}\left[\left(\gamma^{\rho} p_{\rho}-m_{e}\right) \gamma^{\mu}\left(\gamma^{\sigma} p_{\sigma}^{\prime}+m_{e}\right) \gamma^{\nu}\right] \\
&= 4 p_{\rho} p_{\sigma}^{\prime}\left(\eta^{\rho \mu} \eta^{\sigma \nu}-\eta^{\rho \sigma} \eta^{\mu \nu}+\eta^{\rho \nu} \eta^{\mu \sigma}\right)  \tag{5.17}\\
&-m_{e}^{2} 4 \eta^{\mu \nu}+\operatorname{tr}\left(\text { odd number of } \gamma^{\prime} \mathrm{s}\right) \times m_{e} \times \ldots \\
&= 4\left[p^{\mu} p^{\prime \nu}+p^{\nu} p^{\prime \mu}-\left(p \cdot p^{\prime}+m_{e}^{2}\right) \eta^{\mu \nu}\right]
\end{align*}
$$

and similarly

$$
\begin{equation*}
\operatorname{tr}\left[\left(\gamma \cdot k^{\prime}+m_{\mu}\right) \gamma_{\mu}(\gamma \cdot k-m \mu) \gamma_{v}\right]=4\left[k_{\mu}^{\prime} k_{v}+k_{v}^{\prime} k_{\mu}-\left(k \cdot k^{\prime}+m_{\mu}^{2}\right) \eta_{\mu v}\right] . \tag{5.18}
\end{equation*}
$$

Since the ratio of $m_{\mu}$ and $m_{e}$ is about 200 we can drop $m_{e}$. Then altogether after a few cancellations

$$
\begin{equation*}
\frac{1}{4} \sum_{\text {Spins }}|\mathcal{M}|^{2}=8 \frac{e^{4}}{q^{4}}\left[(p \cdot k)\left(p^{\prime} \cdot k^{\prime}\right)+\left(p^{\prime} \cdot k\right)\left(p \cdot k^{\prime}\right)+\left(p \cdot p^{\prime}\right) m_{\mu}^{2}\right] . \tag{5.19}
\end{equation*}
$$

### 5.1.4 Centre-of-mass frame

Switch to the c.o.m. frame and rotate the coordinate frame such that $p$ points in $\hat{z}$-direction and $p^{\prime}$ in - $\hat{\text { z}}$-direction (see Figure 5.2).


Figure 5.2: Centre-of-mass frame, $p$ in $\hat{z}$-direction.

Introducing the angle $\theta$ between $p^{\prime}$ and $k$ and taking the relativistic limit, i.e. $m_{e}^{2} \ll|\vec{p}|^{2}$, yields

$$
\begin{equation*}
p=\binom{E}{E \hat{z}}, \quad p^{\prime}=\binom{E}{-E \hat{z}}, k=\binom{E}{\vec{k}}, \quad \vec{k} \cdot \hat{z}=|\vec{k}| \cos \theta . \tag{5.20}
\end{equation*}
$$

Therefore

$$
\begin{align*}
q^{2} & =\left(p+p^{\prime}\right)^{2}=2 p \cdot p^{\prime}=4 E^{2}, \\
p \cdot k & =E^{2}-E|\vec{k}| \cos \theta=p^{\prime} \cdot k^{\prime},  \tag{5.21}\\
p \cdot k^{\prime} & =E^{2}+E|\vec{k}| \cos \theta=p^{\prime} \cdot k .
\end{align*}
$$

Plugging this in and eliminating $|\vec{k}|^{2}$ using $|\vec{k}|^{2}=E^{2}-m_{\mu}^{2}$ yields

$$
\begin{equation*}
\frac{1}{4} \sum_{\text {Spins }}|\mathcal{M}|^{2}=e^{4}\left[\left(1+\frac{m_{\mu}^{2}}{E^{2}}\right)+\left(1-\frac{m_{\mu}^{2}}{E^{2}}\right) \cos ^{2} \theta\right] . \tag{5.22}
\end{equation*}
$$

### 5.1.5 Cross-section

To this end recall the general formula for a 2-2 scattering event. For this we had derived the expression

$$
\begin{equation*}
\mathrm{d} \sigma=\frac{(2 \pi)^{4}}{4 E_{p} E_{p^{\prime}}\left|v_{p}-v_{p^{\prime}}\right|} \mathrm{d} \Pi_{2} \delta^{(4)}\left(p+p^{\prime}-k-k^{\prime}\right)\left|\mathcal{M}_{f i}\right|^{2} \tag{5.23}
\end{equation*}
$$

with

$$
\begin{equation*}
\mathrm{d} \Pi_{2}=\frac{\mathrm{d}^{3} k}{(2 \pi)^{3}} \frac{\mathrm{~d}^{3} k^{\prime}}{(2 \pi)^{3}} \frac{1}{2 E_{k} 2 E_{k^{\prime}}} \tag{5.24}
\end{equation*}
$$

Note that no factor of $\frac{1}{2}$ is required here since $\mu^{+}$and $\mu^{-}$are distinguishable. In the c.o.m. frame we generally have

$$
\begin{equation*}
\mathrm{d} \Pi_{2}(2 \pi)^{4} \delta^{(4)}\left(p+p^{\prime}-k-k^{\prime}\right)=\frac{\left.\mathrm{d}|\vec{k}| \vec{k}\right|^{2} \mathrm{~d} \Omega}{(2 \pi)^{3} 4 E_{k} E_{k^{\prime}}} \delta\left(E_{\mathrm{com}}-E_{k}-E_{k^{\prime}}\right) \tag{5.25}
\end{equation*}
$$

with $E_{k}=\sqrt{|\vec{k}|^{2}+m_{k}^{2}}$ and $E_{k^{\prime}}=\sqrt{|\vec{k}|^{2}+m_{k^{\prime}}^{2}}$ and

$$
\begin{equation*}
\mathrm{d}|\vec{k}| \delta\left(E_{\mathrm{com}}-E_{k}-E_{k^{\prime}}\right)=\left(\frac{|\vec{k}|}{E_{k}}+\frac{|\vec{k}|}{E_{k^{\prime}}}\right)^{-1} \tag{5.26}
\end{equation*}
$$

Then, altogether

$$
\begin{equation*}
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega}=\frac{1}{4 E_{p} E_{p^{\prime}}\left|v_{p}-v_{p^{\prime}}\right|} \frac{|\vec{k}|}{16 \pi^{2} E_{\mathrm{com}}}\left|\mathcal{M}_{f i}\right|^{2} . \tag{5.27}
\end{equation*}
$$

Note that, if all 4 particles had equal masses and the outstates were indistinguishable, this would correctly reduce to the famous expression encountered earlier

$$
\begin{equation*}
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega}=\frac{1}{2!} \frac{1}{64 \pi^{2}} \frac{1}{s}\left|\mathcal{M}_{f i}\right|^{2} \tag{5.28}
\end{equation*}
$$

with $s=E_{\mathrm{com}}^{2}$. Applied to the present case (5.27) yields, with

$$
\begin{equation*}
E_{p}=E_{p^{\prime}}=E=\frac{1}{2} E_{\mathrm{com}}, \quad\left|v_{p}-v_{p^{\prime}}\right|=2 \tag{5.29}
\end{equation*}
$$

the differential cross-section

$$
\begin{equation*}
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega}=\frac{1}{32 \pi^{2}} \frac{1}{E_{\mathrm{com}}^{2}} \frac{|\vec{k}|}{E_{\mathrm{com}}} \frac{1}{4} \sum_{\text {Spins }}\left|\mathcal{M}_{f i}\right|^{2} \tag{5.30}
\end{equation*}
$$

Introducing the fine-structure constant

$$
\begin{equation*}
\alpha:=\frac{e^{2}}{4 \pi} \sim \frac{1}{137} \tag{5.31}
\end{equation*}
$$

leaves us with

$$
\begin{equation*}
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega}=\frac{\alpha^{2}}{4 E_{\mathrm{com}}^{2}} \sqrt{1-\frac{m_{\mu}^{2}}{E^{2}}}\left[\left(1+\frac{m_{\mu}^{2}}{E^{2}}\right)+\left(1-\frac{m_{\mu}^{2}}{E^{2}}\right) \cos ^{2} \theta\right] \tag{5.32}
\end{equation*}
$$

for the differential cross-section. Thus for the total cross-section

$$
\begin{equation*}
\sigma=\frac{4 \pi \alpha^{2}}{3 E_{\mathrm{com}}} \sqrt{1+\frac{m_{\mu}}{E^{2}}}\left(1+\frac{1}{2} \frac{m_{\mu}^{2}}{E^{2}}\right) . \tag{5.33}
\end{equation*}
$$

The first term reflects the purely kinematic and thus universal energy dependence, while the second term represents the correction due to specifics of the QED interaction and is thus characteristic of the concrete dynamics involved in this process.

### 5.2 The Ward-Takahashi identity

Consider a QED amplitude $\mathcal{M}(k)$ involving an external photon of momentum $k^{\mu}$ with $k^{2}=0$ and polarisation $\xi^{\mu}(k)$. We thus can write the scattering amplitude as

$$
\begin{equation*}
\mathcal{M}(k)=\xi^{\mu}(k) \mathcal{M}_{\mu}(k) . \tag{5.34}
\end{equation*}
$$

Then the Ward-Takahashi identity for QED is the statement that

$$
\begin{equation*}
k^{\mu} \mathcal{M}_{\mu}(k)=0 . \tag{5.35}
\end{equation*}
$$

To appreciate its significance recall from Gupta-Bleuler quantisation that the constraint

$$
\begin{equation*}
\partial \cdot A^{+}\left|\psi_{\text {phys }}\right\rangle=0 \tag{5.36}
\end{equation*}
$$

implies $\zeta^{\mu} k_{\mu}=0$ for states $|k, \zeta\rangle \in \mathcal{H}_{\text {phys }}$. This left 2 positive norm transverse polarisations $\zeta_{T}$ and 1 zero norm polarisation $\zeta_{s}=k$. (5.35) proves that the unphysical zero-norm polarization state $\left|k, \zeta_{s}\right\rangle$ decouples from the $S$-matrix as an external state, as claimed.
(5.35) follows from a more general from of the Ward identity. Consider a theory with arbitrary spin fields $\phi_{a}(x)$ with a global continuous symmetry

$$
\begin{equation*}
\phi_{a}(x) \rightarrow \phi_{a}(x)+\epsilon \delta \phi_{a}(x), \quad \epsilon \in \mathbb{R}, \tag{5.37}
\end{equation*}
$$

and a conserved Noether current $j^{\mu}(x)$ such that classically

$$
\begin{equation*}
\partial_{\mu} j^{\mu}=0 \tag{5.38}
\end{equation*}
$$

holds on-shell. Then in the quantum theory the general Ward-Takahashi-identity is a statement about current conservation inside a general n-point function:

$$
\begin{align*}
0= & \partial_{\mu}\langle\Omega| T j^{\mu}(x) \phi_{a_{1}}\left(x_{1}\right) \ldots \phi_{a_{n}}\left(x_{n}\right)|\Omega\rangle \\
& +i \sum_{j=1}^{n}\langle\Omega| \phi_{a_{1}} \ldots \hat{\phi}_{a_{j}}\left(x_{j}\right) \delta \phi_{a_{j}}\left(x_{j}\right) \delta^{(4)}\left(x-x_{j}\right) \ldots \phi_{a_{n}}\left(x_{n}\right)|\Omega\rangle, \tag{5.39}
\end{align*}
$$

where $\hat{\phi}_{a_{j}}\left(x_{j}\right)$ is omitted. (5.39) can be proven as follows:

- For definiteness and simplicity consider a scalar theory. For the free theory, the classical equation of motion is

$$
\begin{equation*}
\frac{\delta S}{\delta \phi(x)}:=\frac{\partial \mathcal{L}}{\partial \phi(x)}-\partial_{\mu} \frac{\partial L}{\partial\left(\partial_{\mu} \phi(x)\right)}=-\left(\partial^{2}+m^{2}\right) \phi(x)=0 . \tag{5.40}
\end{equation*}
$$

This equation is satisfied by the Heisenberg quantum fields as an operator equation.

- With the help of this operator equation, one finds that

$$
\begin{equation*}
\left(\partial_{x}^{2}+m^{2}\right) i\langle 0| T \phi(x) \phi\left(x_{1}\right)|0\rangle=\delta^{(4)}\left(x-x_{1}\right) \tag{5.41}
\end{equation*}
$$

for a 2-point-function. The $\delta$-distribution occurs due to the action of $\partial_{x^{0}}$ on $\Theta\left(x-x_{0}\right)$ appearing in the time-ordering prescription symbolized by $T$. More generally one finds

$$
\begin{align*}
& \left(\partial_{x}^{2}+m^{2}\right) i\langle 0| T \phi(x) \phi\left(x_{1}\right) \ldots \phi\left(x_{n}\right)|0\rangle \\
& \quad=\sum_{j=1}^{n}\langle 0| T \phi\left(x_{1}\right) \ldots \hat{\phi}\left(x_{j}\right) \delta^{(4)}\left(x-x_{j}\right) \ldots \phi\left(x_{n}\right)|0\rangle . \tag{5.42}
\end{align*}
$$

This can be written as the Schwinger-Dyson equation

$$
\begin{equation*}
\langle 0| T \frac{\delta S}{\delta \phi(x)} \phi\left(x_{1}\right) \ldots \phi\left(x_{n}\right)|0\rangle=i \sum_{j=1}^{n}\langle 0| T \phi\left(x_{1}\right) \ldots \hat{\phi}\left(x_{j}\right) \delta^{(4)}\left(x-x_{j}\right) \ldots \phi\left(x_{n}\right)|0\rangle . \tag{5.43}
\end{equation*}
$$

- By Noether's theorem, if $\phi(x) \rightarrow \phi(x)+\epsilon \delta \phi(x)$ is a global continuous symmetry, then the Noether current enjoys

$$
\begin{equation*}
\partial_{\mu} j^{\mu}=-\left(\frac{\partial \mathcal{L}}{\partial \phi(x)}-\partial_{\mu} \frac{\partial \mathcal{L}}{\partial\left(\partial_{\mu} \phi(x)\right)}\right) \delta \phi(x) \equiv-\frac{\delta S}{\delta \phi(x)} \delta \phi(x) \tag{5.44}
\end{equation*}
$$

off-shell. Plugging this into the Schwinger-Dyson equation yields:

$$
\begin{align*}
& \partial_{\mu}\langle 0| T j^{\mu}(x) \phi\left(x_{1}\right) \ldots \phi\left(x_{n}\right)|0\rangle \\
& \quad+i \sum_{j=1}^{n}\langle 0| T \phi\left(x_{1}\right) \ldots \hat{\phi}\left(x_{j}\right) \delta \phi(x) \delta^{(4)}\left(x-x_{j}\right) \ldots \phi\left(x_{n}\right)|0\rangle=0 . \tag{5.45}
\end{align*}
$$

- In interacting theories the reasoning goes through - all that changes is that the equations of motion involve extra polynomial terms due to the interactions, which however do not alter the conclusions. This yields the corresponding statements about the full correlation functions $\langle\Omega| \ldots|\Omega\rangle$.

The Schwinger-Dyson equation and the Ward-identity show that the classical equation of motion and current conservation hold inside the correlation functions only up to so-called contact terms, which are precisely the extra terms we pick up if the insertion point $x$ of the operator $\frac{\delta S}{\delta \phi(x)}$ or $\partial_{\mu} j^{\mu}(x)$ coincides with the insertion point $x_{j}$ of one of the other fields inside the correlator.

## Caveat:

It is important to be aware that the $n$-point correlator $\langle\Omega| \ldots|\Omega\rangle$ involves in general divergent loopdiagrams to be discussed soon. The Ward identity only holds if the regularisation required to define these divergent intergrals respects the classical symmetry. We will see that for QED such regulators can be found. On the other hand, if no regulator exists that respects the Ward-identity for a classical
symmetry, this symmetry is anomalous - it does not hold at the quantum level. Such anomalies will be studied in great detail in QFT II.

It remains to deduce (5.35) from (5.39):

- Following LSZ, the scattering amplitude $\mathcal{M}(k)$ involving an external photon $\left|k^{\mu}, \xi^{\mu}(k)\right\rangle$ and $n$ further particles is given by

$$
\begin{equation*}
\langle f \mid i\rangle=i Z_{A}^{-1 / 2} \xi^{\mu}(k) \int \mathrm{d}^{4} x e^{-i k \cdot x} \partial_{x}^{2} \ldots\langle\Omega| T A_{\mu}(x) \underbrace{\ldots}_{(n \text { other fields })}|\Omega\rangle \tag{5.46}
\end{equation*}
$$

In Feynman gauge the full interacting QED Lagrangian reads

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{2} \partial_{\mu} A_{\nu} \partial^{\mu} A^{\nu}-A_{\mu} j^{\mu}+\bar{\psi}\left(i \gamma \cdot \partial-m_{0}\right) \psi, \quad j^{\mu}=e \bar{\psi} \gamma^{\mu} \psi \tag{5.47}
\end{equation*}
$$

The classical equation of motion for $A^{\mu}$ is

$$
\begin{equation*}
\partial^{2} A^{\mu}=j^{\mu} \tag{5.48}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\partial_{x}^{2}\langle\Omega| T A_{\mu}(x) \ldots|\Omega\rangle=\langle\Omega| T j_{\mu}(x) \ldots|\Omega\rangle+\text { contact terms }, \tag{5.49}
\end{equation*}
$$

where the contact terms include $(n-1)$ fields inside $\langle\Omega| \ldots|\Omega\rangle$. Being correlators only of $(n-1)$ fields, these contact terms cannot have precisely $n$ poles in the momenta of the $n$ fields and thus do not contribute to $\langle f \mid i\rangle$ according to the LSZ formalism.

- Now we are left with

$$
\begin{equation*}
\langle f \mid i\rangle=i Z_{A}^{-1 / 2} \xi^{\mu} \int \mathrm{d}^{4} x e^{-i k \cdot x}\langle\Omega| T j_{\mu}(x) \ldots|\Omega\rangle \tag{5.50}
\end{equation*}
$$

For $\xi^{\mu}=k^{\mu}$ we find

$$
\begin{align*}
k^{\mu} \int \mathrm{d}^{4} x e^{-i k \cdot x}\langle\Omega| T j_{\mu}(x) \ldots|\Omega\rangle & =i \int \mathrm{~d}^{4} x\left(\partial^{\mu} e^{-i k \cdot x}\langle\Omega| T j_{\mu}(x) \ldots|\Omega\rangle\right. \\
& =-i \int \mathrm{~d}^{4} x e^{-i k \cdot x} \partial^{\mu}\langle\Omega| T j_{\mu}(x) \ldots|\Omega\rangle \tag{5.51}
\end{align*}
$$

where we used that surface terms vanish (because, as in LSZ, we are really having suitable wave-packets in mind). According to (5.39) this is

$$
\begin{equation*}
0+\text { contact terms. } \tag{5.52}
\end{equation*}
$$

Also in this case the contact terms do not contribute to $\langle f \mid i\rangle$ because we are essentially trading one $A_{\mu}(x)$ field against one matter field $\delta \varphi(x)$ inside the correlator so that the resulting pole structure is not of the right form, in and (5.35) is proven.

### 5.2.1 Relation between current conservation and gauge invariance

The proof of $k^{\mu} \mathcal{M}_{\mu}(k)=0$ only requires that $A_{\mu}(x)$ couples to a conserved current in the sense that

$$
\begin{equation*}
S=S_{A}^{0}[A]+S_{\mathrm{int}}[A, \phi]+S_{\text {matter }}^{\mathrm{rest}}[\phi] \tag{5.53}
\end{equation*}
$$

with

$$
\begin{equation*}
\frac{\delta S_{\mathrm{int}}[A, \phi]}{\delta A_{\mu}}=-j^{\mu} \tag{5.54}
\end{equation*}
$$

the Noether current associated with a global continuous symmetry of the full interacting action $S$.

- For a massless vector theory, which, as a free theory, must always be gauge invariant, (5.54) is in fact equivalent to invariance of the theory under combined gauge transformations of the vector and the matter sector.

To prove that (5.54) implies gauge invariance, our point of departure is the existence of a global continuous symmetry $\phi(x) \rightarrow \phi(x)+\epsilon \delta \phi(x)$ with $\epsilon \in \mathbb{R}$ constant, which leaves the full action $S$ invariant. If we perform instead a local transformation $\phi(x) \rightarrow \phi(x)+\epsilon(x) \delta \phi(x)$ with varying $\epsilon(x)$, then $\delta S$ is in general no longer zero. Since it vanishes for constant $\epsilon, \delta S$ must be proportional (to first order) to $\partial_{\mu} \epsilon(x)$. In fact, Lorentz invariance implies that there exists a 4 -vector $\tilde{j}^{\mu}(x)$ such that $\delta S=\int \tilde{j}^{\mu} \partial_{\mu} \epsilon(x)=-\int\left(\partial_{\mu} \tilde{j}^{\mu}\right) \epsilon(x)$. This identifies $\tilde{j}^{\mu}$ with $j^{\mu}$, the conserved current (because for $\epsilon$ constant, the Noether current has the property that $-\int\left(\partial_{\mu} j^{\mu}\right) \epsilon=\int\left(\frac{\delta S}{\delta \phi} \delta \phi \epsilon\right) \equiv \delta S$ - see equ. (5.44)). Now, since by assumption $\frac{\delta S_{\text {int }}}{\delta A_{\mu}}=-j^{\mu}$ we have $\int\left(\frac{\delta S_{\text {int }}}{\delta A_{\mu}}+j^{\mu}\right) \partial_{\mu} \epsilon(x)=0$. Writing $S=S_{A}^{0}[A]+S_{\text {int }}[A, \phi]+S_{\text {matter }}^{\text {rest }}[\phi]$, the fact that $S_{A}^{0}[A]$ is gauge invariant for a massless vector theory implies that this is just the variation of $S$ with respect to the combined gauge transformation

$$
\begin{equation*}
A_{\mu} \rightarrow A_{\mu}+\partial_{\mu} \epsilon(x), \quad \phi(x) \rightarrow \phi(x)+\epsilon(x) \delta \phi(x) . \tag{5.55}
\end{equation*}
$$

The theory is thus gauge invariant.

To prove the other direction, suppose (5.55) leaves the action $S_{A}^{0}[A]+S_{\text {int }}[A, \phi]+S_{\text {matter }}^{\text {rest }}[\phi]$ invariant (where again $S_{A}^{0}[A]$ is separately gauge invariant). Then there exists some conserved $j^{\mu}(x)$ such that $\int\left(\frac{\delta S_{\text {int }}}{\delta A_{\mu}}+j^{\mu}\right) \partial_{\mu} \epsilon(x)=0$, and we conclude (5.54).
To conclude, massless vector theories are consistent if and only if the vector couples to a conserved current. This is the necessary and consistent condition for gauge invariance and for the Ward identities.

- If we couple a massive vector, whose action is never gauge invariant, to a conserved current as in (5.54), then $k^{\mu} \mathcal{M}_{\mu}=0$ still holds. The fact that in a massive vector theory the Ward identities are still satisfied provided that theory couples as in (5.54) is crucial for its consistency: Recall that in a massive vector theory $\left(k^{2}>0\right)$, the negative norm states are the states

$$
\begin{equation*}
\left|k^{\mu}, \xi^{\mu}\right\rangle \text { with } \xi^{\mu}=k^{\mu} \tag{5.56}
\end{equation*}
$$

Restriction to $\zeta^{\mu}$ with $\zeta^{\mu} k_{\mu}=0$ removes these. This is only justified in the interacting theory as long as no dangerous states $\left|k^{\mu}, \zeta^{\mu}=k^{\mu}\right\rangle$ are produced, as is guaranteed thanks to $k^{\mu} \mathcal{M}_{\mu}=0$ for couplings to conserved currents.

### 5.2.2 Photon polarisation sums in QED

To evaluate the consequences of the Ward identities we can take, without loss of generality, the photon 4 -momentum to be

$$
k^{\mu}=k\left(\begin{array}{l}
1  \tag{5.57}\\
0 \\
0 \\
1
\end{array}\right)
$$

and consider the corresponding basis of polarisation vectors

$$
\epsilon^{\mu}(k, 0)=\left(\begin{array}{l}
1  \tag{5.58}\\
0 \\
0 \\
0
\end{array}\right), \quad \epsilon^{\mu}(k, 1)=\left(\begin{array}{l}
0 \\
1 \\
0 \\
0
\end{array}\right), \quad \epsilon^{\mu}(k, 2)=\left(\begin{array}{l}
0 \\
0 \\
1 \\
0
\end{array}\right), \quad \epsilon^{\mu}(k, 3)=\left(\begin{array}{l}
0 \\
0 \\
0 \\
1
\end{array}\right) .
$$

The Ward identity $k^{\mu} \mathcal{M}_{\mu}(k)=0$ then implies $\mathcal{M}^{0}(k)=\mathcal{M}^{3}(k)$. In physical applications we often need to sum over the two transverse polarisation vectors $\lambda=1,2$ of an external photon involved in a scattering process. By the Ward identity this becomes

$$
\begin{align*}
\sum_{\lambda=1}^{2}\left|\epsilon_{\mu}(k, \lambda) \mathcal{M}^{\mu}(k)\right|^{2} & =\sum_{\lambda=1}^{2} \epsilon_{\mu}(k, \lambda) \epsilon_{v}^{*}(k, \lambda) \mathcal{M}^{\mu}(k)\left(\mathcal{M}^{v}(k)\right)^{*}=\left|\mathcal{M}^{1}(k)\right|^{2}+\left|\mathcal{M}^{2}(k)\right|^{2} \\
& =\left|\mathcal{M}^{1}(k)\right|^{2}+\left|\mathcal{M}^{2}(k)\right|^{2}+\left|\mathcal{M}^{3}(k)\right|^{2}-\left|\mathcal{M}^{0}(k)\right|^{2} \\
& =-\eta_{\mu v} \mathcal{M}^{\mu}(k)\left(\mathcal{M}^{v}(k)\right)^{*} \tag{5.59}
\end{align*}
$$

Thus in sums over polarisations of the above type we can replace $\sum_{\lambda=1}^{2} \epsilon_{\mu}(k, \lambda) \epsilon_{\nu}^{*}(k, \lambda)$ by $-\eta_{\mu v}$. This will be used heavily in deriving the Klein-Nishina formula for Compton scattering in the tutorials.

### 5.2.3 Decoupling of potential ghosts

So far we have merely shown that the zero-norm states with polarization $k^{\mu}$ decouple from the Smatrix and are thus not produced in scattering experiments. But what about the status of timelike polarisation states, which correspond to the even more dangerous negative norm ghosts? Even if we declare $\mathcal{H}_{\text {phys }}$ not to contain them, we must prove that no such states are created as outgoing states from physical in-states in scattering processes. Otherwise the interactions would render the theory inconsistent. ${ }^{2}$

[^27]In fact, it is again the Ward identity that guarantees that no ghosts are created as external states in interactions. The argument relies on the relation

$$
\begin{equation*}
\sum_{\lambda=1}^{2} \epsilon_{\mu}(k, \lambda) \epsilon_{\nu}^{*}(k, \lambda) \mathcal{M}^{\mu}(k)\left(\mathcal{M}^{\nu}(k)\right)^{*}=-\eta_{\mu \nu} \mathcal{M}^{\mu}(k)\left(\mathcal{M}^{\nu}(k)\right)^{*} \tag{5.60}
\end{equation*}
$$

derived above and is sketched as follows: It suffices to show that the restriction of the S-matrix $S$ to the space of transverse polarisations is unitary. This means that we lose no information by considering only transverse in and out-states and thus guarantees that no unphysical states can be produced out of transverse incoming states. Let $P$ denote the projector onto these states of transverse polarisations. The relation (5.60) amounts to the statement

$$
\begin{equation*}
S^{\dagger} P S=S^{\dagger} S \tag{5.61}
\end{equation*}
$$

Since $S^{\dagger} S=\mathbb{1}$, this implies that $(P S P)^{\dagger}(P S P)=P$. This is precisely the statement that the restriction of $S$ to the state of transverse polarisation is unitary.

### 5.3 Radiative corrections in QED - Overview

We now enter a quantitative discussion of radiative corrections in Quantum Field Theory as exemplified by loop corrections to QED. As we will see it is sufficient to study the following types of loop corrections modifying the QED building blocks.

## Corrections to the fermion propagator

We recall that the Feynman propagator of the free fermion field is

$$
\begin{equation*}
S_{F}(x-y)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} e^{-i p(x-y)} \underbrace{\frac{\gamma \cdot p+m_{0}}{p^{2}-m_{0}^{2}+i \epsilon}}_{\equiv \frac{\gamma \cdot p-m_{0}+i \epsilon}{i}} . \tag{5.62}
\end{equation*}
$$

Taking into account corrections due to QED interactions we have instead

$$
\begin{align*}
& \langle\Omega| T \psi(x) \bar{\psi}(y)|\Omega\rangle=y \longrightarrow x+\xrightarrow[\longrightarrow]{\Omega}+\underset{\longrightarrow}{\Omega}+\ldots  \tag{5.63}\\
& \equiv y \text { - } x \text {. }
\end{align*}
$$

Let

$$
\begin{equation*}
A-1 \mathrm{IPI}-B \equiv-i \Sigma(p)_{A B} \tag{5.64}
\end{equation*}
$$

denote the amputated 1PI diagram. By Dyson resummation the full propagator then takes the form

$$
\begin{equation*}
y-x=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} e^{-i p(x-y)} \frac{i}{\gamma \cdot p-m_{0}-\sum(p)+i \epsilon}, \tag{5.65}
\end{equation*}
$$

where $\sum(p)$ is called self-energy of the electron.

## Corrections to the photon propagator

The Fourier transform of the photon propagator, denoted by $\sim \sim \sim$, derives by Dyson resummation from the 1-PI diagram

$$
\begin{equation*}
\mu \sim \Omega \text { (1P) } \sim v \equiv i \Pi^{\mu \nu}\left(q^{2}\right)=\text { self-energy of the photon or vacuum polarisation. } \tag{5.66}
\end{equation*}
$$

## Corrections to the interaction vertex

Loop corrections modifiy the cubic interaction vertex. We will find it useful to define an effective vertex by summing up all loop-corrections. Schematically,


We will compute these corrections to 1 -loop order. The diagrams will exhibit

- ultraviolet (UV) divergences from integrating the momenta of particles in the loop up to infinity and
- infra-red (IR) divergences if the diagram contains massless particles - i.e. photons - running in the loop.

The general status of these singularities is as follows:

- IR divergences in loop-diagrams cancel against IR divergences from radiation of soft, i.e. lowenergy photons and thus pose no conceptual problem.
- UV divergences require regularisation of the integral and can be absorbed in a clever definition of the parameters via renormalisation.


### 5.4 Self-energy of the electron at 1-loop

At 1-loop order the electron propagator takes the form

$$
\begin{equation*}
\langle\Omega| T \psi(x) \bar{\psi}(y)|\Omega\rangle=y \longrightarrow x+\underbrace{\substack{\Omega}}_{=\int \frac{d^{4} p}{\frac{d^{4} p}{(2 \pi)^{4}} e^{i p(x-y) \times(\mathrm{I})}}} \tag{5.67}
\end{equation*}
$$

where the integrand (I) is given by the diagram

$$
\begin{equation*}
(\mathrm{I}) \equiv p \xrightarrow{\widetilde{S}_{k}} p . \tag{5.68}
\end{equation*}
$$

The photon in the loop carries momentum $p-k$. By means of the Feynman rules,

$$
\begin{equation*}
(\mathrm{I})=\frac{i\left(\gamma \cdot p+m_{0}\right)}{p^{2}-m_{0}^{2}+i \epsilon}\left(-i \Sigma_{2}(p)\right) \frac{i\left(\gamma \cdot p+m_{0}\right)}{p^{2}-m_{0}^{2}+i \epsilon} . \tag{5.69}
\end{equation*}
$$

The amputated 1-loop contribution corresponds to omitting the two outer fermion propagators and is thus given by

$$
\begin{equation*}
i \Sigma_{2}(p)=(-i e)^{2} \int \frac{\mathrm{~d}^{4} k}{(2 \pi)^{4}} \gamma^{\mu} \frac{i\left(k+m_{o}\right)}{k^{2}-m_{0}^{2}+i \epsilon} \gamma^{\nu} \frac{-i \eta_{\mu \nu}}{(p-k)^{2}+i \epsilon} . \tag{5.70}
\end{equation*}
$$

It will turn out that the integral is divergent near $k=0$ if $p \rightarrow 0$. This is an IR divergence. A careful analysis reveals that it will cancel in all amplitudes against similar such IR divergences from other diagrams involving soft photons and thus poses no harm. For the present discussion we could just ignore it, but for completeness we introduce a ficticious small photon mass $\mu$ to regulate the IR divergence:

$$
\begin{equation*}
i \Sigma_{2}(p)=(-i e)^{2} \int \frac{\mathrm{~d}^{4} k}{(2 \pi)^{4}} \gamma^{\mu} \frac{i\left(k+m_{o}\right)}{k^{2}-m_{0}^{2}+i \epsilon} \gamma_{\mu} \frac{-i}{(p-k)^{2}-\mu^{2}+i \epsilon} . \tag{5.71}
\end{equation*}
$$

The evaluation of such typical momentum integrals proceeds in 3 steps:

### 5.4.1 Feynman parameters

The integrand contains a fraction of the form $\frac{1}{A B}$ with $A=(p-k)^{2}-\mu^{2}+i \epsilon$ and $B=k^{2}-m_{0}^{2}+i \epsilon$. It turns out useful to write this as $\frac{1}{(\ldots)^{2}}$ and to complete the square in $k$. To this end we exploit the elementary identity

$$
\begin{equation*}
\frac{1}{A B}=\int_{0}^{1} \mathrm{~d} x \frac{1}{(x A+(1-x) B)^{2}} \tag{5.72}
\end{equation*}
$$

$x$ is called a Feynman parameter. Applying this identity yields in the present case

$$
\begin{align*}
\frac{1}{A B} & =\int_{0}^{1} \mathrm{~d} x \frac{1}{\left(x\left((p-k)^{2}-\mu^{2}+i \epsilon\right)+(1-x)\left(k^{2}-m_{0}^{2}+i \epsilon\right)\right)^{2}} \\
& =\int_{0}^{1} \mathrm{~d} x \frac{1}{\left(k^{2}-2 x k \cdot p+x p^{2}-x \mu^{2}-(1-x) m_{0}^{2}+i \epsilon+x^{2} p^{2}-x^{2} p^{2}\right)^{2}}  \tag{5.73}\\
& =\int_{0}^{1} \mathrm{~d} x \frac{1}{\left(l^{2}-\Delta+i \epsilon\right)^{2}},
\end{align*}
$$

where $l=k-x p$ and $\Delta=-x(1-x) p^{2}+x \mu^{2}+(1-x) m_{0}^{2}$. Then

$$
\begin{equation*}
-i \Sigma_{2}(p)=-e^{2} \int_{0}^{1} \mathrm{~d} x \int \frac{\mathrm{~d}^{4} l}{(2 \pi)^{4}} \frac{\gamma^{\mu}\left(k+m_{0}\right) \gamma_{\mu}}{\left(l^{2}-\Delta+i \epsilon\right)^{2}} \tag{5.74}
\end{equation*}
$$

The term in the numerator can be simplified with the help of the gamma-matrix identity

$$
\begin{equation*}
\gamma^{\mu} \gamma^{\nu} k_{v} \gamma_{\mu}=\left(2 \eta^{\mu \nu}-\gamma^{v} \gamma^{\mu}\right) k_{v} \gamma_{\mu}=k^{\mu} \gamma_{\mu}\left(2-\gamma^{\nu} \gamma_{v}\right) \tag{5.75}
\end{equation*}
$$

and, in $d$-dimensions, $\gamma^{\mu} \gamma_{\mu}=d$. Therefore altogether (keeping $d$ arbitrary for later purposes)

$$
\begin{equation*}
\gamma^{\mu}\left(k+m_{0}\right) \gamma_{\mu}=(2-d)(l+x p)+d m_{0} . \tag{5.76}
\end{equation*}
$$

Now for symmetry reasons

$$
\begin{equation*}
\int \frac{\mathrm{d}^{4} l}{(2 \pi)^{4}} \frac{l^{\mu}}{\left(l^{2}-\Delta\right)^{2}}=0 \tag{5.77}
\end{equation*}
$$

and thus, if $d \equiv 4$,

$$
\begin{equation*}
-i \Sigma_{2}(p)=-e^{2} \int_{0}^{1} \mathrm{~d} x \int \frac{\mathrm{~d}^{4} l}{(2 \pi)^{4}} \frac{-2 x p+4 m_{0}}{\left(l^{2}-\Delta+i \epsilon\right)^{2}} \tag{5.78}
\end{equation*}
$$

Remark: For more general loop integrals one makes use of the identity

$$
\begin{equation*}
\frac{1}{A_{1} \ldots A_{n}}=\int_{0}^{1} \mathrm{~d} x_{1} \ldots \mathrm{~d} x_{n} \delta\left(\sum_{i} x_{i}-l\right) \frac{(n-1)!}{\left(x_{1} A_{1}+\ldots+x_{n} A_{n}\right)^{n}} \tag{5.79}
\end{equation*}
$$

which can be proven by induction.

### 5.4.2 Wick rotation

As we have seen, we encounter loop-integrals of the typical form $\int \frac{\mathrm{d}^{4} l}{(2 \pi)^{4}} \frac{1}{\left(l^{2}-\Delta+i \epsilon\right)^{n}}$. This integral would be relatively easy to perform if it were defined in Euclidean space. The Wick rotation relates it to such a Euclidean integral.
Indeed, as the $i \epsilon$ factors remind us, the $l^{0}$ integral is in fact a complex contour integral along the real axis. The value of this integral is unchanged if we deform the contour without hitting any pole. Therefore we can rotate the contour by $90^{\circ}$ counter-clockwise to lie along the imaginary axis. Introducing the Euclidean 4-momentum $l_{E}=\left(l_{E}^{0}, \vec{l}_{E}\right)$ as

$$
\begin{equation*}
l^{0}=i l_{E}^{0}, \quad \vec{l}=\vec{l}_{E} \tag{5.80}
\end{equation*}
$$

such that $l^{2} \equiv-l_{E}^{2} \equiv-\sum_{i}\left(l_{E}^{i}\right)^{2}$, we can write the integral as

$$
\begin{equation*}
\int \frac{\mathrm{d}^{4} l}{(2 \pi)^{4}} \frac{1}{\left(l^{2}-\Delta+i \epsilon\right)^{m}}=i(-1)^{m} \int \frac{\mathrm{~d}^{4} l_{E}}{(2 \pi)^{4}} \frac{1}{\left(l_{E}^{2}+\Delta-i \epsilon\right)^{m}} \tag{5.81}
\end{equation*}
$$

Since we won't need the $i \epsilon$ any longer we can omit it at this stage. We can now peform the integral as a spherical integral in $\mathbb{R}^{4}$.

### 5.4.3 Regularisation of the integral

The integral

$$
\begin{equation*}
\mathrm{I}_{4}=\int \frac{\mathrm{d}^{4} l_{E}}{(2 \pi)^{4}} \frac{1}{\left(l_{E}^{2}+\Delta\right)^{2}}=\int_{0}^{\infty} \frac{\left|l_{E}\right|^{3} \mathrm{~d}\left|l_{E}\right|}{(2 \pi)^{4}} \int \mathrm{~d} \Omega_{4} \frac{1}{\left(\left|l_{E}\right|^{2}+\Delta\right)^{2}} \tag{5.82}
\end{equation*}
$$

is divergent due to integration over the UV region $\left|l_{E}\right| \rightarrow \infty$. To isolate the divergence we regularise the integral. The 3 most common methods for regularisation are

- Momentum cutoff: Isolate the UV divergence as a divergence in the upper momentum limit $\Lambda$ by writing

$$
\begin{equation*}
\int_{0}^{\infty} \ldots=\lim _{\Lambda \rightarrow \infty} \int_{0}^{\Lambda} \ldots \tag{5.83}
\end{equation*}
$$

Following this procedure we can evaluate the integral in an elementary fashion as

$$
\begin{align*}
\mathrm{I}_{4} & =\int \mathrm{d} \Omega_{4} \lim _{\Lambda \rightarrow \infty} \int_{0}^{\infty} \frac{\mathrm{d}\left|l_{E}\right|}{(2 \pi)^{4}} \frac{\left|l_{E}\right|^{3}}{\left(\left|l_{E}\right|^{2}+\Delta\right)^{2}} \\
& =\int \mathrm{d} \Omega_{4} \lim _{\Lambda \rightarrow \infty} \frac{1}{2} \int_{0}^{\infty} \frac{\mathrm{d}\left|l_{E}\right|^{2}}{(2 \pi)^{4}} \frac{\left|l_{E}\right|^{2}}{\left(\left|l_{E}\right|^{2}+\Delta\right)^{2}}  \tag{5.84}\\
& =\int \mathrm{d} \Omega_{4} \lim _{\Lambda \rightarrow \infty} \frac{1}{2} \frac{1}{(2 \pi)^{4}}\left[\frac{\Delta}{\Delta+\left|l_{E}\right|^{2}}+\log \left(\left|l_{E}\right|^{2}+\Delta\right)\right]_{R^{2}=0}^{\Lambda^{2}} \\
& =\int \mathrm{d} \Omega_{4} \lim _{\Lambda \rightarrow \infty} \frac{1}{2} \frac{1}{(2 \pi)^{4}}\left[\log \left(\frac{\Lambda^{2}}{\Delta}\right)-1\right] .
\end{align*}
$$

The loop leads to a log-divergence as $\Lambda \rightarrow \infty$. The problem is that this regularisation procedure is not consistent with the Ward identities, as we will see when computing the photon propagator, and is therefore not a useful regularisation method in QED.

- Dimensional regularisation (dimReg) is probably the most common method. We first evaluate the integral in $d$ dimensions; writing $d=4-\epsilon$ then isolates the divergence as a pole in $\epsilon$ as $\epsilon \rightarrow 0$.

Let us therefore consider

$$
\begin{equation*}
\mathrm{I}_{d}=\frac{\mathrm{d}^{d} l_{E}}{(2 \pi)^{d}} \frac{1}{\left(l_{E}^{2}+\Delta\right)^{2}}=\int \frac{\mathrm{d} \Omega_{d}}{(2 \pi)^{3}} \int_{0}^{\infty} \mathrm{d}\left|l_{E}\right| \frac{\left|l_{E}\right|^{d-1}}{\left(\left|l_{E}\right|^{2}+\Delta\right)^{2}} \tag{5.85}
\end{equation*}
$$

The volume of the unit sphere in $d$-dimensions is

$$
\begin{equation*}
\int \mathrm{d} \Omega_{d}=\frac{2 \pi^{d / 2}}{\Gamma\left(\frac{d}{2}\right)} \tag{5.86}
\end{equation*}
$$

where

$$
\begin{equation*}
\Gamma(z)=\int_{0}^{\infty} \mathrm{d} y y^{z-1} e^{-y} \tag{5.87}
\end{equation*}
$$

is the Euler $\Gamma$-function. This is because

$$
\begin{align*}
(\sqrt{\pi})^{d} & =\left(\int_{-\infty}^{\infty} \mathrm{d} x e^{-x^{2}}\right)^{d}=\int \mathrm{d}^{d} x e^{-\sum_{i=1}^{d} x_{i}^{2}} \\
& =\int \mathrm{d} \Omega_{d} \int_{0}^{\infty} \mathrm{d}|x \| x|^{d-1} e^{-|x|^{2}}  \tag{5.88}\\
& =\int \mathrm{d} \Omega_{d} \frac{1}{2} \int_{0}^{\infty} \mathrm{d}\left(x^{2}\right)\left(x^{2}\right)^{d / 2-1} e^{-x^{2}} .
\end{align*}
$$

The $\Gamma$-function has the properties

- $\Gamma(n)=(n-1)$ ! for $n \in \mathbb{N}^{+}$as can be shown by integration by parts.
- $\Gamma(z)$ has analytic poles at $z=0,-1,-2,-3, \ldots$; for a proof we refer to standard textbooks on complex analysis.

Continuing with our integral we write

$$
\begin{equation*}
\int_{0}^{\infty} \mathrm{d}\left|l_{E}\right| \frac{\left|l_{E}\right|^{d-1}}{\left(\left|l_{E}\right|^{2}+\Delta\right)^{2}}=\frac{1}{2} \int_{0}^{\infty} \mathrm{d}\left|l_{E}\right|^{2} \frac{\left(\left|l_{E}\right|^{2}\right)^{d / 2-1}}{\left(\left|l_{E}\right|^{2}+\Delta\right)^{2}} \tag{5.89}
\end{equation*}
$$

and substitute

$$
\begin{equation*}
x=\frac{\Delta}{\left|l_{E}\right|^{2}+\Delta} \Rightarrow\left|l_{E}\right|^{2}=\frac{\Delta}{x}-\Delta \tag{5.90}
\end{equation*}
$$

with

$$
\begin{equation*}
\mathrm{d} x=-\mathrm{d}\left|l_{E}\right|^{2} \frac{\Delta}{\left(\left|l_{E}\right|^{2}+\Delta\right)^{2}} . \tag{5.91}
\end{equation*}
$$

Then the integral (5.89) becomes

$$
\begin{equation*}
=\frac{1}{2}\left(\frac{1}{\Delta}\right)^{2-d / 2} \int_{0}^{1} \mathrm{~d} x x^{1-d / 2}(1-x)^{d / 2-1} . \tag{5.92}
\end{equation*}
$$

As a last ingredient we need the Euler $\beta$-function

$$
\begin{equation*}
\mathcal{B}(\alpha, \beta):=\int_{0}^{1} \mathrm{~d} x x^{\alpha-1}(1-x)^{\beta-1}=\frac{\Gamma(\alpha) \Gamma(\beta)}{\Gamma(\alpha+\beta)}, \tag{5.93}
\end{equation*}
$$

where the last identity is non-trivial and proven again in textbooks on complex analysis.
We can rewrite (5.92) in terms of $\mathcal{B}(\alpha, \beta)$ with $\alpha=2-d / 2$ and $\beta=d / 2$. Together with $\Gamma(2)=1$ we arrive at

$$
\begin{equation*}
\frac{1}{2}\left(\frac{1}{\Delta}\right)^{2-d / 2} \int_{0}^{1} \mathrm{~d} x x^{1-d / 2}(1-x)^{d / 2-1}=\frac{1}{2}\left(\frac{1}{\Delta}\right)^{2-d / 2} \Gamma(2-d / 2) \Gamma(d / 2) \tag{5.94}
\end{equation*}
$$

Therefore

$$
\begin{equation*}
\mathrm{I}_{d}=\frac{\mathrm{d}^{d} l_{E}}{(2 \pi)^{d}} \frac{1}{\left(l_{E}^{2}+\Delta\right)^{2}}=\frac{1}{(4 \pi)^{d / 2}} \Gamma(2-d / 2)\left(\frac{1}{\Delta}\right)^{2-d / 2} \tag{5.95}
\end{equation*}
$$

We now apply this to $\mathrm{I}_{d}$ with $d=4$ and note that we encounter a singularity from $\Gamma(0)$. To isolate the singularity we define

$$
\begin{equation*}
d=4-\epsilon \tag{5.96}
\end{equation*}
$$

and would like to write $\mathrm{I}_{4} \propto \lim _{\epsilon \rightarrow 0} \mathrm{I}_{4-\epsilon}$. However $\mathrm{I}_{4}$ is dimensionless because

$$
\begin{equation*}
\left[l_{E}\right]=[\operatorname{mass}] \tag{5.97}
\end{equation*}
$$

whereas $\mathrm{I}_{4-\epsilon}$ has dimension [mass] ${ }^{-\epsilon}$. If we introduce a compensating mass scale $\tilde{M}$, we can write

$$
\begin{align*}
\mathrm{I}_{4} & =\lim _{\epsilon \rightarrow 0} \tilde{M}^{\epsilon} \mathrm{I}_{4-\epsilon}=\lim _{\epsilon \rightarrow \infty} \tilde{M}^{\epsilon} \frac{1}{4 \pi^{2-\epsilon / 2}} \Gamma(\epsilon / 2)\left(\frac{1}{\Delta}\right)^{\epsilon / 2}  \tag{5.98}\\
& =\lim _{\epsilon \rightarrow 0} \frac{1}{4 \pi^{2}} \Gamma(\epsilon / 2)\left(\frac{4 \pi \tilde{M}^{2}}{\Delta}\right)^{\epsilon / 2}
\end{align*}
$$

Finally we use
$-\Gamma(\epsilon / 2)=\frac{2}{\epsilon}-\gamma+O(\epsilon)$, where $\gamma \approx 0.5772$ is the Euler-Mascheroni number and we refer again to the complex analysis literature for a proof;

- $x^{\epsilon / 2}=1+\frac{\epsilon}{2} \log (x)+O\left(\epsilon^{2}\right)$, which is just a Taylor expansion in $\epsilon$ since $x^{\epsilon / 2}=e^{\frac{\epsilon}{2} \log (x)}$.

With this input

$$
\begin{equation*}
\mathrm{I}_{4}=\lim _{\epsilon \rightarrow 0} \frac{1}{(4 \pi)^{2}}\left(\frac{2}{\epsilon}-\gamma+\log \left(\frac{4 \pi \tilde{M}^{2}}{\Delta}\right)+O(\epsilon)\right) \tag{5.99}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\mathrm{I}_{4}=\lim _{\epsilon \rightarrow 0} \frac{1}{(4 \pi)^{2}}\left(\frac{2}{\epsilon}+\log \left(\frac{M^{2}}{\Delta}\right)+O(\epsilon)\right) \tag{5.100}
\end{equation*}
$$

where $M^{2}=4 \pi e^{-\gamma} \tilde{M}^{2}$.
The final result for the electron self-energy at 1-loop is

$$
\begin{equation*}
-i \Sigma_{2}(p)=\left.(-i) \frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x\left((2-\epsilon / 2) m_{0}-(1-\epsilon / 2) x \not p\right)\left(\frac{2}{\epsilon}+\log \frac{M^{2}}{\Delta}\right)\right|_{\epsilon \rightarrow 0} \tag{5.101}
\end{equation*}
$$

where $M$ is some as yet arbitrary mass-scale. We will understand its significance soon.

- Pauli-Villars (PV) regularisation: We subtract a diagram with a ficiticious massive particle in the loop - here a photon - of mass $\Lambda$, i.e. we compute


Subtracting the second diagram subtracts the divergence because for infinite momenta the mass of the particle in the loop becomes irrelevant and both diagrams asymptote to the same divergent value. If $\Lambda \rightarrow \infty$ the contribution of the subtracted vanishes and we recover the actual diagram. The divergence reappears therefore as a divergence in $\Lambda$ as $\Lambda \rightarrow \infty$.

Subtracting both diagrams yields

$$
\begin{align*}
-i \Sigma_{2}(p) & =\lim _{\Lambda \rightarrow \infty}(\text { computation with } \mu-\text { computation with } \Lambda) \\
& =\ldots \\
-i \Sigma_{2}(p) & =\lim _{\Lambda \rightarrow \infty}(-i) \frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x\left(2 m_{0}-x p\right) \log \left(\frac{x \Lambda^{2}}{\Delta}\right), \tag{5.103}
\end{align*}
$$

where the omitted steps are elementary and similar to the integration performed in (5.84). We can rewrite this by introducing again an arbitrary mass scale $M$ as

$$
\begin{equation*}
-i \Sigma_{2}(p)=\left.\lim _{\Lambda \rightarrow \infty}(-i) \frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x\left(2 m_{0}-x p\right)\left(\log \left(\frac{x \Lambda^{2}}{M^{2}}\right)+\log \left(\frac{M^{2}}{\Delta}\right)\right)\right|_{\Lambda \rightarrow \infty} \tag{5.104}
\end{equation*}
$$

Comparing (5.101) and (5.104) we note that the momentum-dependent terms are in both regularisation approaches the same, while the $\frac{1}{\epsilon}$-divergence in dimReg corresponds to a logarithmic divergence in PV or in momentum cutoff regularisation.

### 5.5 Bare mass $m_{0}$ versus physical mass $m$

Before dealing with the UV divergence of the electron-propagator, let us recall the non-perturbative information encoded in the 2-point function. By the Källén-Lehmann spectral representation we have

$$
\begin{equation*}
\int \mathrm{d}^{4} x e^{i p \cdot x}\langle\Omega| T \psi(x) \bar{\psi}(0)|\Omega\rangle=\frac{i}{p-m_{0}-\Sigma(p)} \stackrel{!}{=} \frac{i Z_{2}}{p-m}+\text { terms analytic at } m . \tag{5.105}
\end{equation*}
$$

The physical mass $m$ is the location of the lowest-lying analytic pole and can thus be computed by solving

$$
\begin{equation*}
p-m_{0}-\left.\Sigma(p)\right|_{p=m}=0, \tag{5.106}
\end{equation*}
$$

while the wavefunction renormalisation of the electron - called $Z_{2}$ - is the residue of the propagator at $p=m$. To find this residue we perform a Taylor expansion of

$$
\begin{equation*}
f(p):=p-m_{0} \mathbb{1}-\Sigma(p) \tag{5.107}
\end{equation*}
$$

around $p_{0}=m \cdot \mathbb{1}$. Then

$$
\begin{align*}
f(p) & =f\left(p_{0}\right)+\left.\left(p-p_{0}\right) \frac{\mathrm{d} f}{\mathrm{~d} p}\right|_{p_{0}}+O\left(\left(p-p_{0}\right)^{2}\right) \\
& =0+\left.(p-m \cdot \mathbb{1})\left(1-\frac{\mathrm{d} \Sigma(p)}{\mathrm{d} p}\right)\right|_{p=m \cdot \mathbb{1}}+O\left((p-m)^{2}\right) . \tag{5.108}
\end{align*}
$$

Therefore

$$
\begin{equation*}
\left.\operatorname{res}\left[\frac{1}{p-m_{0}-\Sigma(p)}\right]\right|_{p=m}=\left.\left(1-\frac{\mathrm{d} \Sigma(p)}{\mathrm{d} p}\right)^{-1}\right|_{p=m}=: Z_{2} \tag{5.109}
\end{equation*}
$$

and thus

$$
\begin{equation*}
Z_{2}^{-1}=1-\left.\frac{\mathrm{d} \Sigma(p)}{\mathrm{d} p}\right|_{p=m} \tag{5.110}
\end{equation*}
$$

We can now compute $m$ to order $\alpha$, where

$$
\begin{equation*}
\Sigma(p)=\Sigma_{2}(p)+O\left(\alpha^{2}\right) \tag{5.111}
\end{equation*}
$$

Then to order $\alpha$

$$
\begin{equation*}
m-m_{0}-\Sigma_{2}(p=m)=0 \tag{5.112}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\delta m:=m-m_{o}=\Sigma_{2}(p=m)+O\left(\alpha^{2}\right) . \tag{5.113}
\end{equation*}
$$

Note that if we trade $m$ by $m_{0}$ in $\Sigma_{2}(p=m)$, the error will be $O\left(\alpha^{2}\right)$ and thus we can also write

$$
\begin{equation*}
\delta m:=m-m_{o}=\Sigma_{2}\left(p=m_{0}\right)+O\left(\alpha^{2}\right) . \tag{5.114}
\end{equation*}
$$

In dimReg we find

$$
\begin{equation*}
m-m_{0}=\frac{\alpha}{2 \pi} m_{0} \int_{0}^{1} \mathrm{~d} x\left((2-x)+\frac{\epsilon}{2}(x-1)\right)\left[\frac{2}{\epsilon}+\log \left(\frac{M^{2}}{(1-x)^{2} m_{0}^{2}+x \mu^{2}}\right)\right] \tag{5.115}
\end{equation*}
$$

and in PV

$$
\begin{equation*}
m-m_{0}=\frac{\alpha}{2 \pi} m_{0} \int_{0}^{1} \mathrm{~d} x(2-x) \underbrace{\left[\log \left(\frac{x \Lambda^{2}}{M^{2}}\right)+\log \left(\frac{M^{2}}{(1-x)^{2} m_{0}^{2}+x \mu^{2}}\right)\right]}_{\equiv \log \left(\frac{x \Lambda^{2}}{(1-x)^{2} m_{0}^{2}+x \mu^{2}}\right)} . \tag{5.116}
\end{equation*}
$$

Note that $m-m_{0}>0$ : This is because what we are computing is literally the self-energy of the electron, i.e. the mass shift due to the (positive!) energy stored in its own electric field. Recall from classical electrodynamics that also classically this quantity is divergent due to the pointlike structure of the electron. Not surprisingly, the divergence remains in QED. ${ }^{3}$
Likewise we can compute $Z_{2}$ to order $\alpha$ by evaluating (5.110) at order $\alpha$, the result being

$$
\begin{equation*}
Z_{2}=1+\frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x\left[-x \log \left(\frac{x \Lambda^{2}}{(1-x)^{2} m^{2}+x \mu^{2}}\right)+2(2-x) \frac{x(1-x) m^{2}}{(1-x)^{2} m^{2}+x \mu^{2}}\right]+O\left(\alpha^{2}\right) \tag{5.117}
\end{equation*}
$$

[^28]
### 5.5.1 Mass renormalisation

We thus encounter an obvious problem: $m-m_{0}$ is divergent due to the UV divergence of the propagator. What saves the day are the following crucial observations:

- Only $m$, the physical mass, is a physical observable. Namely $m$ is the rest mass of an electron as measured in experiments.
- By contrast, $m_{0}$, the so-called bare mass, is merely a parameter that appears in the Lagrangian and per se cannot be measured directly. Rather the Lagrangian produces for us, via the Feynman rules, measurable quantities, the scattering amplitudes, which depend on $m_{0}$.

This suggests the following solution to the divergence of $\delta m$ : The divergence can be absorbed in the definition of $m_{0}$ by interpreting the equation (5.113) for $\delta m$ as an equation for $m_{0}$ in terms of the measured physical mass $m$ and the cutoff $\Lambda$ or $\epsilon$. Concretely in PV regularisation (to be specific)

$$
\begin{equation*}
m_{0}=m\left(1-\frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x(2-x) \log \left(\frac{x \Lambda^{2}}{(1-x)^{2} m_{0}^{2}+x \mu^{2}}\right)\right)+\boldsymbol{O}\left(\alpha^{2}\right)=m_{0}(\Lambda) . \tag{5.118}
\end{equation*}
$$

This means that we take the parameter $m_{0}$ in the Lagrangian to be divergent. We compute scattering amplitudes etc. in terms of this divergent object $m_{0}(\Lambda)$ and at the end plug in the above equation for $m_{0}(\Lambda)$ to express everything via $m$. If through this procedure all physical quantities in the end are independent of $\Lambda$, the theory is said to be renormalisable. Indeed QED is renormalisable.We will discuss how one can see this more systematically later.
Rather than being systematic here, let us exemplify this in the following trivial example: Compute the mass $m$ at 1-loop. Well, this means we take the equation for $m$ in terms of $m_{0}$ at 1-loop and plug in (5.118). Obviously we find $m=m$. This demonstrates that absorbing the divergence in the bare mass comes at a price: We lose predicitivity for the physical mass. In other words, we must now take $m$ directly from experiment. In a renormalisable theory we retain, however, predicitivity for all but a finite number of quantities.

### 5.6 The photon propagator

The building block to compute radiative corrections to the photon propagator is the 1PI amputated diagram

$$
\begin{equation*}
v \sim \sim_{q} \simeq \mathbb{1 P I} \sim \sim_{q} \sim \mu:=i \Pi^{\mu \nu}(q) . \tag{5.119}
\end{equation*}
$$

On general grounds we can make the following statements:

- By Lorentz invariance its tensorial structure can only depend on $\eta^{\mu \nu}$ and $q^{\mu} q^{\nu}$.
- The Ward identity further implies transversality,

$$
\begin{equation*}
q^{\mu} \Pi_{\mu \nu}(q)=0, \tag{5.120}
\end{equation*}
$$

because we can view $\Pi^{\mu \nu}$ as a $1-1$ scattering amplitude. ${ }^{4}$ Thus

$$
\begin{equation*}
i \Pi^{\mu v}(q)=\left(\eta^{\mu v}-\frac{q^{\mu} q^{v}}{q^{2}}\right) f\left(q^{2}\right) \tag{5.121}
\end{equation*}
$$

Note in particular that the Ward identity implies that no terms of the form $\frac{q^{\mu} q^{\nu}}{m_{0}^{2}}$ or the alike arise as these would destroy transversality.

- Finally, $i \Pi^{\mu \nu}(q)$ cannot have an analyic pole at $q^{2}=0,{ }^{5}$ because this would require a singleparticle massless intermediate state (whose propagator vanishes at zero momentum), but no such intermediate states occur for the 1 PI diagram relevant for the photon propagator.


Figure 5.3: Possible diagrams: The first loop carries only massive particles. In the second diagram the photon in the loop is massless, but it does not arise as a single-particle due to the accompanying electron line.

Therefore,

$$
\begin{equation*}
i \Pi^{\mu \nu}(q)=\left(q^{2} \eta^{\mu \nu}-q^{\mu} q^{v}\right) \Pi\left(q^{2}\right) \tag{5.122}
\end{equation*}
$$

such that $\Pi\left(q^{2}\right)$ is regular at $q^{2}=0$. It is useful to define the projection operator onto momenta orthogonal to $q^{\mu}$,

$$
\begin{equation*}
P^{\mu v}(q)=\eta^{\mu \nu}-\frac{q^{\mu} q^{v}}{q^{2}} \tag{5.123}
\end{equation*}
$$

with

$$
\begin{equation*}
q_{\mu} P^{\mu \nu}(q)=0 \text { and } P^{\mu \nu}(q) P_{\nu \rho}(q)=P_{\rho}^{\mu}(q) \tag{5.124}
\end{equation*}
$$

By Dyson resummation the Fourier transform of the full propagator takes the form

$$
\begin{align*}
& =\frac{-i}{q^{2}\left(1-\Pi\left(q^{2}\right)\right)}\left(\eta_{\mu \nu}-\frac{q_{\mu} q_{\nu}}{q^{2}}\right)+\frac{-i}{q^{2}} \frac{q_{\mu} q_{\nu}}{q^{2}}, \tag{5.125}
\end{align*}
$$

where the final result follows with the help of (5.124) as will be discussed in the tutorial. Note that this holds in Feynman gauge $\xi=1$. In general gauge we would get ${ }^{6}$

$$
\begin{equation*}
\sim \sim \sim=\frac{-i}{q^{2}\left(1-\Pi\left(q^{2}\right)\right)}\left(\eta^{\mu v}-\frac{q^{\mu} q^{v}}{q^{2}}\right)+\xi \frac{-i}{q^{2}} \frac{q^{\mu} q^{v}}{q^{2}} \tag{5.126}
\end{equation*}
$$

[^29]This identifies the $\xi$-dependent term as pure gauge. We can omit it by going to Landau gauge $\xi=0$, in which

$$
\begin{equation*}
\left.\sim \sim \sim\right|_{\text {Landau }}=\frac{-i}{q^{2}\left(1-\Pi\left(q^{2}\right)\right)}\left(\eta_{\mu \nu}-\frac{q_{\mu} q_{v}}{q^{2}}\right) \tag{5.127}
\end{equation*}
$$

This result has two important consequences:

- Since $\Pi\left(q^{2}\right)$ is regular at $q^{2}=0$ (again in the sense that for finite cutoff, there is no analytic pole) the pole at $q^{2}=0$ is unaffected by the radiative corrections. Thus the photon remains massless as required by gauge invariance. This is a consequence of the Ward identity, which had lead to this form of $i \Pi^{\mu \nu}(q)$.
- The photon field strength renormalisation is simply

$$
\begin{equation*}
Z_{3}=\frac{1}{1-\Pi(0)} . \tag{5.128}
\end{equation*}
$$

To obtain the result at order $\alpha$ we procced in a similar manner as for $\Sigma_{2}(p)$. We denote by $\Pi_{2}\left(q^{2}\right)$ the 1-loop contribution to $\Pi\left(q^{2}\right)$ (just like $\Sigma_{2}(p)$ denotes the 1-loop contribution to $\Sigma(p)$ ). To compute $\Pi_{2}\left(q^{2}\right)$ we must consider the diagram

and bring the result into the form (5.121) with $\Pi_{2}\left(q^{2}\right)$ instead of $\Pi\left(q^{2}\right)$.
This computation is most conveniently carried out in dimensional regularization. We merely quote the final result

$$
\begin{equation*}
\Pi_{2}\left(q^{2}\right)=-\left.\frac{2 \alpha}{\pi} \int_{0}^{1} \mathrm{~d} x x(1-x)\left(\frac{2}{\epsilon}+\log \frac{M^{2}}{\Delta}\right)\right|_{\epsilon \rightarrow 0}, \tag{5.130}
\end{equation*}
$$

where $\Delta=m_{0}^{2}-x(1-x) q^{2}$.
To conclude this section we state without proof that if we were to compute the 1-loop corrected propagator with a naive momentum cutoff $\int_{0}^{\infty} \mathrm{d} p \rightarrow \int_{0}^{\Lambda} \mathrm{d} p$, then we would find a divergent term of the form

$$
\begin{equation*}
i \Pi_{2}^{\mu \nu}(q) \sim \eta^{\mu \nu} \Lambda^{2}+P^{\mu \nu}(\ldots) . \tag{5.131}
\end{equation*}
$$

This would violate $q^{\mu} \Pi_{\mu \nu}(q)=0$ and thus the Ward identity. This shows, as claimed, that momentum cutoff regularisation breaks gauge invariance and is thus not useful in QED. Note that from our derivation of the Ward identities, it may not be completely obvious why cutoff regularization leads to a breakdown of the Ward identities. Next term we will get to know a very simple derivation of the Ward identities in the path integral quantization approach and interpret the breakdown of the Ward identities as due to the fact that the cutoff-regularised measure of the path integral is not invariant under the $\mathrm{U}(1)$ symmetry of the classical action.

### 5.7 The running coupling

Radiative corrections to the photon propagator are responsible for an important phenomenon in QED, the running of the electric coupling. To see this we consider a scattering process with an intermediate photon of the form


Taking into account the appearance of $e_{0}$ at each vertex, it is clear that the amplitude involves a factor of

$$
\begin{equation*}
\frac{\left(-i e_{0}\right)^{2}}{1-\Pi\left(q^{2}\right)} \frac{-i\left(\eta^{\mu \nu}-q^{\mu} q^{\nu} / q^{2}\right)}{q^{2}} \tag{5.132}
\end{equation*}
$$

We can absorb the correction term $\left(1-\Pi\left(q^{2}\right)\right)^{-1}$ into the coupling and define an effective coupling

$$
\begin{equation*}
e\left(q^{2}\right):=\frac{e_{0}}{\sqrt{1-\Pi\left(q^{2}\right)}} \tag{5.133}
\end{equation*}
$$

This obviously depends on $q^{2}$, the energy transferred by the photon in the scattering process. When measuring the physical charge of an electron, we therefore need to specify the energy scale at which the measurement is performed. Let us define the physical coupling or renormalised charge as the effective charge as measured at $q^{2}=0$, i.e.

$$
\begin{equation*}
e:=\lim _{q^{2} \rightarrow 0} e\left(q^{2}\right)=\frac{e_{0}}{\sqrt{1-\Pi(0)}}=e_{0} \sqrt{Z_{3}} . \tag{5.134}
\end{equation*}
$$

Note that $\Pi(0)$ is a divergent constant. E.g. at 1-loop in perturbation theory it takes the form in dimensional regularisation

$$
\begin{equation*}
\Pi(0)=\Pi_{2}(0)+O\left(\alpha^{2}\right)=-\frac{2 \alpha}{\pi} \int_{0}^{1} \mathrm{~d} x x(1-x)\left(\frac{2}{\epsilon}+\log \frac{M^{2}}{m_{0}^{2}}\right)+O\left(\alpha^{2}\right) \tag{5.135}
\end{equation*}
$$

Note that at this order in $\alpha, m_{0}$ and $m$ can be exchanged (because the difference is itself of $O(\alpha)$ ). We proceed as we did when defining the physical electron mass and absorb the divergence into the definition of the bare coupling $e_{0}$. This can be done because only $e$ is a physical observable, whose finite value we take from experiment. We define the bare coupling as

$$
\begin{equation*}
e_{0}=e \sqrt{1-\Pi(0)} \equiv e_{0}(\Lambda) \text { or } e_{0}(\epsilon), \tag{5.136}
\end{equation*}
$$

depending on the regularisation method we used (Pauli-Villars or dimensional regularisation). The effective coupling at energy $q^{2}$ is therefore

$$
\begin{equation*}
e^{2}\left(q^{2}\right)=\frac{e_{0}^{2}}{1-\Pi\left(q^{2}\right)}=\frac{e_{0}^{2}}{1-\Pi_{2}\left(q^{2}\right)}+O\left(\alpha^{2}\right)=\frac{e^{2}\left(1-\Pi_{2}(0)\right)}{1-\Pi_{2}\left(q^{2}\right)}+O\left(\alpha^{2}\right) \tag{5.137}
\end{equation*}
$$

where $1-\Pi_{2}(0)=\frac{1}{1+\Pi_{2}(0)}+O\left(\alpha^{2}\right)$. Thus

$$
\begin{equation*}
e^{2}\left(q^{2}\right)=\frac{e^{2}}{1-\left(\Pi_{2}\left(q^{2}\right)-\Pi_{2}(0)\right)}+\boldsymbol{O}\left(\alpha^{2}\right) \tag{5.138}
\end{equation*}
$$

where $e^{2} \equiv e^{2}\left(q^{2}=0\right)$ is the physical charge at $q^{2}=0$. We conclude that the effective coupling at $q^{2} \neq 0$ is

$$
\begin{equation*}
e^{2}\left(q^{2}\right)=\frac{e^{2}}{1-\hat{\Pi}_{2}\left(q^{2}\right)}+O\left(\alpha^{2}\right) \tag{5.139}
\end{equation*}
$$

with $\hat{\Pi}_{2}\left(q^{2}\right)=\Pi_{2}\left(q^{2}\right)-\Pi_{2}(0)$ finite and independent of $\Lambda$. Concretely,

$$
\begin{equation*}
\hat{\Pi}_{2}\left(q^{2}\right)=-\frac{2 \alpha}{\pi} \int_{0}^{1} \mathrm{~d} x x(1-x) \log \left(\frac{m^{2}}{m^{2}-x(1-x) q^{2}}\right) \tag{5.140}
\end{equation*}
$$

The effective coupling of QED increases logarithmically with the energy scale at which the experiment is performed. The physical interpretation of this is that the charge of an electron is screened by virtual $e^{+} e^{-}$pairs so the effective charge decreases at long distance corresponding to small energy. Thus the QED vacuum appears like a polarisable medium. This explains the name vacuum polarisation for the radiatively corrected photon propagator.

## A comment on the Landau pole

Note that from this analysis, the effective coupling increases indefinitely as we increase the energy. This phenomenon is called the Landau pole of QED and casts doubt on the validity of the theory at arbitrarily high energies. We will find in QFT II that suitable non-abelian gauge theories such as QCD exhibit precisely the opposite phenomenon: they become asymptotically free and are thus perfectly well-defined in the UV. One proposed solution to the Landau pole problem of QED is therefore that the electromagnetic gauge group, which is part of the Standard Model, might in fact emerge as part of a unified non-abelian gauge group, whose dynamics takes over at high energies.

### 5.8 The resummed QED vertex

We define the amputated resummed cubic QED vertex as


For later purposes we point out that in diagrams with external fermions this resummed vertex is dressed with suitable factors of $Z_{2}$. For instance the amplitude associated with the diagram

is of the form

$$
\begin{equation*}
i \mathcal{M} \sim \bar{u}\left(p^{\prime}\right)\left(-i e_{0} Z_{2} \Gamma^{\mu}\left(p^{\prime}, p\right)\right) u(p) \frac{-i\left(\eta_{\mu \nu}-q^{\mu} q^{\nu} / q^{2}\right)}{q^{2}\left(1-\Pi\left(q^{2}\right)\right.} \bar{u}\left(k^{\prime}\right)\left(-i e_{0} Z_{2} \Gamma^{\nu}\left(k^{\prime}, k\right)\right) u(k) . \tag{5.141}
\end{equation*}
$$

Perturbatively the cubic vertex can be expanded as


By Lorentz invariance and the Ward identity, which implies that

$$
\begin{equation*}
q^{\mu} \bar{u}\left(p^{\prime}\right) \Gamma_{\mu}\left(p^{\prime}, p\right) u(p)=0 \tag{5.142}
\end{equation*}
$$

one can show that it must take the form

$$
\begin{equation*}
\Gamma^{\mu}\left(p^{\prime}, p\right)=\gamma^{\mu} F_{1}\left(q^{2}\right)+i \frac{S^{\mu \nu} q_{v}}{2 m_{0}} F_{2}\left(q^{2}\right) \tag{5.143}
\end{equation*}
$$

where $F_{1}\left(q^{2}\right)$ and $F_{2}\left(q^{2}\right)$ are called form factors and $S^{\mu \nu}=\frac{i}{4}\left[\gamma^{\mu}, \gamma^{\nu}\right]$. The form factors can be computed perturbatively via loop-integrals. We merely quote the following results:

- $F_{1}\left(q^{2}\right)$ is given by:

$$
\begin{equation*}
F_{1}\left(q^{2}\right)=1+\delta F_{1}\left(q^{2}\right)+O\left(\alpha^{2}\right), \tag{5.144}
\end{equation*}
$$

where $\delta F_{1}\left(q^{2}\right)$ is a UV divergent function, which, e.g. in PV regularization, is given by

$$
\begin{align*}
\delta F_{1}\left(q^{2}\right)= & \frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x \mathrm{~d} y \mathrm{~d} z \delta(x+y+z-1) \times  \tag{5.145}\\
& \times\left[\log \frac{z \Lambda^{2}}{\Delta}+\frac{1}{\Delta}(1-x)(1-y) q^{2}+\left(1-4 z+z^{2}\right) m_{0}^{2}\right],
\end{align*}
$$

with $\Delta=-x y q^{2}+(1-z)^{2} m_{0}^{2}+\mu^{2} z$. Here $\mu$ shows up to regulate the IR divergences that will cancel eventually.

- For $F_{2}\left(q^{2}\right)=0+\delta F_{2}\left(q^{2}\right)+O\left(\alpha^{2}\right)$ one finds in PV regularisation

$$
\begin{equation*}
\delta F_{2}\left(q^{2}\right)=\frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x \mathrm{~d} y \mathrm{~d} z \delta(x+y+z-1) \frac{2 m_{0}^{2} z(1-z)}{\Delta} \tag{5.146}
\end{equation*}
$$

which is finite. This finiteness of $F_{2}\left(q^{2}\right)$ persists to all orders in perturbation theory.
In the limit $p^{\prime} \rightarrow p$, i.e. $q \rightarrow 0$, we expect the loop corrections to vanish, i.e. as far as the vertex corrections are concerned e.g.


Indeed one can show that

$$
\begin{equation*}
\lim _{q \rightarrow 0} Z_{2} \Gamma^{\mu}(p+q, p)=\gamma^{\mu} . \tag{5.147}
\end{equation*}
$$

Note the appearance of the factor $Z_{2}$ due to the external fermion legs, as pointed out before. We can prove (5.147) in two ways:

- By direct inspection in perturbation theory one finds

$$
\begin{align*}
Z_{2} \Gamma^{\mu}(p+q, p) & =\left(1+\delta Z_{2}+O\left(\alpha^{2}\right)\right)\left(\gamma^{\mu}\left(1+\delta F_{1}\left(q^{2}+O\left(\alpha^{2}\right)\right)+i \frac{S^{\mu v} q_{v}}{2 m_{0}} F_{2}\left(q^{2}\right)\right)\right.  \tag{5.148}\\
& =\left(1+\delta Z_{2}+\delta F_{1}\left(q^{2}\right)+O\left(\alpha^{2}\right)\right) \gamma^{\mu}+\frac{i}{2 m_{0}} S^{\mu \nu} q_{v} F_{2}\left(q^{2}\right),
\end{align*}
$$

where $\delta F_{1}\left(q^{2}\right)$ and $\delta Z_{2}$ are related as

$$
\begin{equation*}
\delta F_{1}(0)=-\delta Z_{2} \tag{5.149}
\end{equation*}
$$

This implies (5.147) because $F_{2}\left(q^{2}\right)$ is finite.

- In fact one can prove (5.147) non-perturbatively via the Ward identities. One defines a quantity $Z_{1}$ by

$$
\begin{equation*}
\lim _{q \rightarrow 0} \Gamma^{\mu}(p+q, p)=Z_{1}^{-1} \gamma^{\mu} \tag{5.150}
\end{equation*}
$$

As we will prove on Assignment 12, the Ward identities imply that $Z_{1}=Z_{2}$ and thus

$$
\begin{equation*}
\lim _{q \rightarrow 0} Z_{2} \Gamma^{\mu}(p+q, p)=\gamma^{\mu} \tag{5.151}
\end{equation*}
$$

In turn this proves that the relation (5.149) must hold order by order in perturbation theory.

### 5.8.1 Physical charge revisited

Since the cubic vertex contains information about the coupling strength, one might wonder whether the structure of radiative corrections included in $\Gamma^{\mu}\left(p^{\prime}, p\right)$ is consistent with our previous definition of the effective coupling given in section (5.7) based only on the photon propagator.
To investigate this consider a typical fully resummed diagram of the form


$$
\sim \frac{-i e_{0} \Gamma^{\mu}\left(p^{\prime}, p\right) Z_{2}}{\sqrt{1-\Pi\left(q^{2}\right)}} \frac{1}{q^{2}} \frac{-i e_{0} \Gamma_{\mu}\left(k^{\prime}, k\right) Z_{2}}{\sqrt{1-\Pi\left(q^{2}\right)}}
$$

As $q^{2} \rightarrow 0$ we find $\Gamma^{\mu}(p+q, p) \rightarrow \gamma^{\mu} Z_{1}^{-1}$ (see eq. (5.150)) and thus as $q^{2} \rightarrow 0$ the amplitude reduces to (ignoring polarizations of external fields)

$$
\begin{equation*}
\frac{-i e_{0}}{\sqrt{1-\Pi(0)}} \frac{Z_{2}}{Z_{1}} \gamma^{\mu} \frac{1}{q^{2}} \gamma_{\mu} \frac{-i e_{0}}{\sqrt{1-\Pi(0)}} \frac{Z_{2}}{Z_{1}} . \tag{5.152}
\end{equation*}
$$

Since $Z_{2} / Z_{1}=1$ by the Ward identites, the physical charge at $q^{2}=0$ is $e_{0} \sqrt{Z_{3}}$ as we found before.

### 5.8.2 Anomalous magnetic moment

It is very instructive to study the the non-relativistic limit of the radiatively corrected vertex and compare the interactions it induces with quantum mechanical scattering amplitudes in Born's approximation. This analysis is performed in detail e.g. in Peskin-Schröder, p.187/188. The computation shows, amongst other things, that the coupling of an electron to an external $\vec{B}$-field is described, in the non-relatvistic limit, by a potential

$$
\begin{equation*}
V(\vec{x})=-\vec{\mu} \cdot \vec{B}(\vec{x}), \quad \vec{\mu}=g \frac{e}{2 m} \vec{S}, \tag{5.153}
\end{equation*}
$$

where $\vec{\mu}$ represents the magnetic moment of the electron and $\vec{S}$ denotes the quantum mechanical spin operator. The Landé factor comes out as

$$
\begin{equation*}
g=2\left(\left(F_{1}(0)+Z_{2}-1\right)+F_{2}(0)\right)=2+\underbrace{2 F_{2}(0)}_{=O(\alpha)} . \tag{5.154}
\end{equation*}
$$

The value $g=2$ follows already from relativistic Quantum Mechanics, in which the Dirac equation is interpreted as an equation for the wavefunction of the electron. Crucially, $F_{2}(0)$ yields QED loop corrections to $g=2$. These can be computed order by order in perturbation theory and are in impressive agreement with experiment.

### 5.9 Renormalised perturbation theory of QED

For a systematic treatment of UV divergent diagrams, it suffices to consider all amputated, 1PI UVdivergent diagrams. All divergent diagrams in a QFT are given either by these diagrams or possibly by diagrams containing these as subdiagrams.
In QFT II we will find a simple way to classify the divergent 1PI amputated diagrams in a given QFT. Applied to QED, this classification will prove that in QED the UV divergent 1PI amputated diagrams are precisely the three types of diagrams which we have studied in the previous sections:


Let us recap their properties.

- At 1-loop order we have found the following structure of UV divergences:

$$
\begin{equation*}
i \Pi^{\mu \nu}\left(q^{2}\right)=i\left(\eta^{\mu \nu}-\frac{q^{\mu} q^{\nu}}{q^{2}}\right) \Pi_{2}\left(q^{2}\right) \tag{5.156}
\end{equation*}
$$

Expanding $\Pi_{2}\left(q^{2}\right)$ as a Taylor series in $q^{2}$ yields

$$
\begin{equation*}
\Pi_{2}\left(q^{2}\right)=c_{0}^{(1)} \log \frac{\Lambda}{M}+\text { finite } \times O\left(q^{2}\right) \tag{5.157}
\end{equation*}
$$

That is, the UV divergence appears at order $\left(q^{2}\right)^{0}$ and is characterized by a constant coefficient $c_{0}^{(1)}$ independent of $q^{2}$.

The remaining two diagrams at 1-loop order have the following structure:

and


Therefore the UV divergences are specified, at 1-loop order, by altogether 4 divergent constants.

- In QFT II we will argue that this structure persists to all orders in perturbation theory, i.e.

$$
\begin{align*}
-i \Sigma(p) & =a_{0} m_{0} \log \frac{\Lambda}{M}+a_{1} p \log \frac{\Lambda}{M}+\text { finite } \times O\left(p^{2}\right) \\
a_{0} & =\underbrace{a_{0}^{(1)}}_{O(\alpha)}+\underbrace{a_{0}^{(2)}}_{O\left(\alpha^{2}\right)}+\ldots \quad a_{1}=a_{1}^{(1)}+a_{1}^{(2)}+\ldots \\
\Pi\left(q^{2}\right) & =c_{0} \log \frac{\Lambda}{M}+\text { finite } \times O\left(q^{2}\right)+\ldots  \tag{5.158}\\
c_{0} & =c_{0}^{(1)}+c_{0}^{(2)}+\ldots \\
i \Gamma_{2}^{\mu}\left(p^{\prime}, p\right) & =b_{0} \gamma^{\mu} \log \frac{\Lambda}{M}+\text { finite } \times O\left(\left(p^{\prime}-p\right)^{2}\right) \\
b_{0} & =b_{0}^{(1)}+b_{0}^{(2)}+\ldots
\end{align*}
$$

Thus at each order in perturbation theory, we encounter 4 constants that multiply UV divergent terms in the above 1PI diagrams. One can absorb these 4 divergent constants order by order in perturbation theory by a procedure called renormalisation.
There are two different, but equivalent ways to perform this procedure, which we now discuss in the context of QED.

### 5.9.1 Bare perturbation theory

So far we have worked in bare perturbation theory, which works as follows:

- We start with the bare Lagrangian

$$
\begin{equation*}
\mathcal{L}=\mathcal{L}_{0}\left(e_{0}, m_{0}\right), \tag{5.159}
\end{equation*}
$$

where $e_{0}$ and $m_{0}$ are the so-called bare charge and mass.

- We compute the above three 1PI amputated UV divergent amplitudes to a given order in perturbation theory as functions of $e_{0}$ and $m_{0}$ and a cutoff $\Lambda$ or $\epsilon$. We keep the cutoff $\Lambda$ finite (or $\epsilon$ non-zero) for the time being so that all computations are perfectly well-defined.
- From these amplitudes we deduce

$$
\begin{align*}
m & =m\left(m_{0}, e_{0}, \Lambda\right) \quad \text { as the physical electron mass, } \\
e & =e\left(m_{0}, e_{0}, \Lambda\right) \quad \text { as the physical coupling at } q^{2}=0 \tag{5.160}
\end{align*}
$$

and also

$$
\begin{align*}
& Z_{2}=Z_{2}\left(m_{0}, e_{0}, \Lambda\right), \quad \text { the wavefunction renormalisation of the electron }  \tag{5.161}\\
& Z_{3}=Z_{3}\left(m_{0}, e_{0}, \Lambda\right), \quad \text { the wavefunction renormalisation of } A^{\mu} .
\end{align*}
$$

- The renormalisation step amounts to interpreting (5.160) as an equation for the bare mass $m_{0}$ and the bare coupling $e_{0}$ in terms of the finite physical quantities $m$ and $e$, i.e. we write

$$
\begin{equation*}
m_{0}=m_{0}(e, m, \Lambda), \quad e_{0}=e_{0}(e, m, \Lambda) . \tag{5.162}
\end{equation*}
$$

Thus 2 linear combinations of the 4 UV divergent constants are now contained in $m_{0}$ and $e_{0}$. We plug the expression for $e_{0}$ and $m_{0}$ back into $\mathcal{L}$. The resulting renormalized Lagrangian is now cuttoff dependent,

$$
\begin{equation*}
\mathcal{L}=\mathcal{L}\left(e_{0}(m, e, \Lambda), m_{0}(m, e, \Lambda)\right) \tag{5.163}
\end{equation*}
$$

and in particular divergent if we take $\Lambda \rightarrow \infty$. This is not a problem because the Lagrangian per se has no physical meaning - only physical observables computed from $\mathcal{L}$ must be finite to make sense.

The remaining 2 linear combinations of divergent constants which are not contained in $m_{0}$ and $e_{0}$ are contained in the expressions for $Z_{2}, Z_{3}$. Even though these do not appear in $\mathcal{L}$ in the present formulation, they enter the computation of scattering amplitudes via the Feynman rules.

- We now compute a given observable from $\mathcal{L}$ in perturbation theory to given order as functions of $e_{0}$ and $m_{0}, Z_{2}, Z_{3}$ and $\Lambda$. The UV divergent terms in $\Lambda$ cancel in all final expressions and all observables are finite expressions of $m$ and $e$ plus terms in $\Lambda$ which vanish as we take $\Lambda \rightarrow 0$. For this to work it is crucial that we compute the observable to the same order in perturbation theory to which we have computed (5.160) and (5.161). E.g. if we evaluate (5.160) and (5.161) at 1 -loop order, then we must compute all remaining scattering amplitudes to 1 -loop order as well.
- At the very end we take $\Lambda \rightarrow \infty$ or $\epsilon \rightarrow 0$. All observables remain finite.

The cancellation of $\Lambda$ works because in all amplitudes we have just the correct numbers of $Z_{2}, Z_{3}$ etc. so that all divergences drop out. Rather than give a general proof we demonstrate this for one of the diagrams contributing to Compton scattering, including radiative corrections:


Ignoring boring polarisation factors, we can organize the amplitude as follows:

$$
\begin{equation*}
\left(Z_{3}^{1 / 2} i e_{0} \Gamma^{\mu} Z_{2}\right) \frac{1}{Z_{2}} \frac{i}{p-m_{0}-\Sigma(p)}\left(Z_{3}^{1 / 2} i e_{0} \Gamma_{\mu} Z_{2}\right) . \tag{5.164}
\end{equation*}
$$

We note the following crucial points:

- $e_{0} Z_{3}^{1 / 2}=e$ is finite.
- $\Gamma^{\mu}(p+q, p) Z_{2}$ is also finite. To see this recall that divergence in $\Gamma^{\mu}(p+q, p)$ arises as the $q^{2}$-independent term as parametrized in (5.158) and is thus already contained in $\lim _{q^{2} \rightarrow 0} \Gamma^{\mu}(p+$ $q, p)$. Since $Z_{2}$ is independent of $q^{2}$ it thus surfices to consider $\lim _{q^{2} \rightarrow 0} \Gamma^{\mu}(p+q, p) Z_{2}$. But as discussed around equ. (5.151) this is finite to all orders in perturbation theory by means of the Ward identities.

Explicitly, this can be confirmed perturbatively from

$$
\begin{equation*}
\Gamma^{\mu}(p+q, p)=\gamma^{\mu} F_{1}\left(q^{2}\right)+i \frac{S^{\mu v} q_{v}}{2 m_{0}} F_{2}\left(q^{2}\right) \tag{5.165}
\end{equation*}
$$

Concerning the second term, $F_{2}\left(q^{2}\right)$ is finite as function of $m_{0}$ and $e_{0}$. To a given order in perturbation theory we can replace $e_{0}$ by $e$ and $m_{0}$ by $m$ as the difference is relevant only at the next order and thus the second term is finite, order by order in perturbation theory. Concerning the first term,

$$
\begin{equation*}
F_{1}\left(q^{2}\right)=1+\delta F_{1}\left(q^{2}\right)+O\left(\alpha^{2}\right) \tag{5.166}
\end{equation*}
$$

with

$$
\begin{equation*}
\delta F_{1}\left(q^{2}\right)=\delta F_{1}(0)+f\left(q^{2}\right) \tag{5.167}
\end{equation*}
$$

where $\delta F_{1}(0)$ carries all UV divergences, while $f\left(q^{2}\right)$ is finite as a function of $m_{0}$. Furthermore we have

$$
\begin{equation*}
Z_{2}=1+\delta Z_{2}+O\left(\alpha^{2}\right) \tag{5.168}
\end{equation*}
$$

and

$$
\begin{equation*}
Z_{2} \Gamma^{\mu}(p+q, p)=\gamma^{\mu}(1+\underbrace{\delta Z_{2}+\delta F_{1}(0)}_{=0}+f\left(q^{2}\right)+O\left(\alpha^{2}\right)+\ldots \text { finite }) . \tag{5.169}
\end{equation*}
$$

Therefore the divergence has cancelled out, because $\delta F_{1}(0)=-\delta Z_{2}$. This persists to all orders.

- What remains is the term

$$
\begin{equation*}
\frac{1}{Z_{2}} \frac{i}{p-m_{0}-\Sigma(p)}=\frac{1}{Z_{2}}\left(\frac{i Z_{2}}{p-m}+\text { terms analytic at } p=m\right) . \tag{5.170}
\end{equation*}
$$

Crucially, the terms analytic at $p=m$ do not contain any divergence in the cutoff. This follows from the Taylor expansion (5.108), according to which these terms are given by the second derivative of $\Sigma(p)$, together with the fact that in $-i \Sigma(p)$ all terms quadratic in $p$ and higher are UV finite. The $Z_{2}$ factors in the first term cancel and the entire expression is finite.

Finally let us outline the generalization and consequences of this renormalisation procedure:

- Consider a general QFT. The theory is called renormalisable if only a finite number of resummed amputated 1PI diagrams is UV divergent.
- Suppose the renormalisable QFT contains $m$ different fields and suppose that the UV divergent 1PI diagrams give rise to $n$ divergent constants order by order in perturbation theory. Then ( $n-$ $m$ ) of these constants can be absorbed in the definition of $(n-m)$ unphysical parameters, the so-called bare couplings. This procedure requires specifying the outcome of $(n-m)$ physical observables as external input. The remaining $m$ constants can be absorbed in the definition of the kinetic terms of the $m$ fields without reducing the predictability of the theory further.
- Thus in a renormalisable theory only a finite number $(n-m)$ of physical observables must be specified order by order in perturbation theory, and predictive power is retained for all remaining observables, which can be computed and are finite as we remove the cutoff.
- However, the price to pay for the appearance of the UV divergences in the first place is that the $(n-m)$ observables cannot be computed by the theory even in principle!For instance, QED cannot make any prediction whatsoever for the absolute value of the electron mass or the charge at $q^{2}=0 .^{7}$ As a result the renormalized QFT necessarily contains free parameters that must be fitted to experiment. Another example for such an observable for which no prediction can be made in a QFT with divergent partition function is the vacuum energy (cosmological constant). From a modern and widely accepted point of view (introduced by K. Wilson), a non-UV finite, but renormalisable QFT is an effective theory: The UV divergences hint at a breakdown of the theory at high energies, where it does not describe the microscopic degrees of freedom correctly. Renormalisation hides our ignorance about the true physics at high energies in the $(n-m)$ observables and we can fit the theory to experiment as one typically does with a phenomenological model.

If we want to go beyond this and describe a truly fundamental (as opposed to effective) theory, we need a theory that is UV finite even before renormalisation. As of this writing the only known theory with this property which also includes gravity is string theory. It has indeed no free parameters.

### 5.9.2 Renormalised Perturbation theory

An equivalent treatment is given by so-called renormalised perturbation theory. The aim is to organise perturbation theory directly in the physical parameters $m$ and $e$ and without the need of including $Z_{2}$ and $Z_{3}$ in the Feynman rules. This comes at the expense of certain divergent counterterms in the renormalised Lagrangian. The systematics is as follows:

- Start again with the bare Lagrangian

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F^{2}+\bar{\psi}\left(i \not \partial-m_{0}\right) \psi-e_{0} \bar{\psi} \gamma^{\mu} \psi A_{\mu} . \tag{5.171}
\end{equation*}
$$

[^30]- Recall that the 2-point functions are normalised such that

$$
\begin{align*}
& \sim \sim=\frac{-i Z_{3}\left(\eta_{\mu v}+\ldots\right)}{q^{2}}+\text { analytic terms, }  \tag{5.172}\\
& \sim=\frac{i Z_{2}}{p-m}+\text { analytic terms. }
\end{align*}
$$

We can absorb $Z_{2}$ and $Z_{3}$ into the fields by renormalising the field strengths as

$$
\begin{align*}
A^{\mu} & =: Z_{3}^{1 / 2} A_{r}^{\mu}, \\
\psi & =: Z_{2}^{1 / 2} \psi_{r} \tag{5.173}
\end{align*}
$$

such that no factor $Z_{2}$ and $Z_{3}$ appears in the propagators of $\psi_{r}$ and $A_{r}^{\mu}$,

$$
\begin{equation*}
\langle\Omega| T A_{r}^{\mu} A_{r}^{\nu}|\Omega\rangle \sim \frac{-i \eta^{\mu \nu}+\ldots}{q^{2}}+\ldots, \quad\langle\Omega| T \psi_{r} \bar{\psi}_{r}|\Omega\rangle \sim \frac{i}{p-m}+\ldots \tag{5.174}
\end{equation*}
$$

We can write the same Lagrangian in terms of the renormalised fields as

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} Z_{3} F_{r}^{2}+Z_{2} \bar{\psi}_{r}\left(i \not \partial-m_{0}\right) \psi_{r}-Z_{2} Z_{3}^{1 / 2} e_{0} \bar{\psi}_{r} \gamma^{\mu} \psi_{r} A_{\mu} . \tag{5.175}
\end{equation*}
$$

To compute an amplitude, we apply the Feynman rules but with $m_{0}$ replaced by $Z_{2} m_{0}$ and $e_{0}$ replaced by $Z_{2} Z_{3}^{1 / 2} e_{0}$, and no factors of $Z_{2}^{1 / 2}$ or $Z_{3}^{1 / 2}$ for external particles in the $S$-matrix. It is important to appreciate that

Renormalising the field strengths does not change any physics if we modify the Feynman rules accordingly.

- We can further rewrite this same Lagrangian as follows:

$$
\begin{align*}
\mathcal{L}= & \left(-\frac{1}{4} F_{r}^{2}+\bar{\psi}_{r}(i \not \partial-m) \psi_{r}-e \bar{\psi}_{r} \gamma^{\mu} \psi_{r} A_{r \mu}\right) \\
& +\left(-\frac{1}{4} \delta_{3} F_{r}^{2}+\bar{\psi}_{r}\left(i \delta_{2} \not \partial-\delta_{m}\right) \psi-e \delta_{1} \bar{\psi}_{r} \gamma^{\mu} \psi_{r} A_{r \mu}\right)  \tag{5.176}\\
\equiv & \mathcal{L}_{r}^{(1)}+\mathcal{L}_{r}^{(2)},
\end{align*}
$$

where

$$
\begin{equation*}
\delta_{3}=Z_{3}-1, \quad \delta_{2}=Z_{2}-1, \quad \delta_{m}=Z_{2} m_{0}-m, \quad \delta_{1}=\frac{e_{0}}{e} Z_{2} Z_{3}^{1 / 2}-1=: Z_{1}-1 \tag{5.177}
\end{equation*}
$$

This is only a rewriting of $\mathcal{L}$ in the form (5.175) by adding and subtracting $\mathcal{L}_{r}^{(1)} . \mathcal{L}_{r}^{(2)}$ contains the so-called counterterms. Note that we have defined $Z_{1}$ by $Z_{1} e=e_{0} Z_{2} Z_{3}^{1 / 2}$. At this stage $m$ and $e$ are just arbitrary parameters to be fixed soon.

- The Feynman rules associated with the form (5.176) of the Lagrangian are now as follows: Associated with $\mathcal{L}_{r}^{(1)}$ are the usual Feynman rules, but with the correct couplings as appearing
in $\mathcal{L}_{r}^{(1)}$ :

$$
\begin{align*}
& \left.\longrightarrow=\frac{-i \eta_{\mu \nu}}{q^{2}+i \epsilon} \quad \text { (in Feynman gauge }\right), \\
& =-i e \gamma^{\mu}, \quad \text { interms of } e, \text { not } e_{0} . \tag{5.178}
\end{align*}
$$

The counterterms in $\mathcal{L}_{r}^{(2)}$ give rise to additional diagrams (counterterm diagrams). Their structure becomes evident if we view the terms in $\mathcal{L}_{r}^{(2)}$ as extra couplings leading to amputated diagrams. In deriving the Feynman rules note that a derivative $\partial^{\mu}$ in position space will give rise to a factor of $-i p^{\mu}$ in the momentum space Feynman rules.

- Thus the counterterm

$$
\begin{equation*}
\bar{\psi}_{r}\left(i \delta_{2} \not \partial-\delta_{m}\right) \psi_{r} \tag{5.179}
\end{equation*}
$$

gives rise to

$$
\begin{equation*}
\longrightarrow \quad \cdots \quad i\left(p \delta_{2}-\delta_{m}\right) . \tag{5.180}
\end{equation*}
$$

- The coupling of a photon to two fermions, i.e. $-e \delta_{1} \psi_{r} \gamma^{\mu} \psi_{r} A_{r \mu}$, gives rise to

$$
\begin{equation*}
-i e \gamma^{\mu} \delta_{1} . \tag{5.181}
\end{equation*}
$$

- The coupling of two photons

$$
\begin{equation*}
-\frac{1}{4} \delta_{3} F_{r \mu \nu} F_{r}^{\mu \nu}=-\frac{1}{2} \delta_{3} A_{r \mu}\left(-\eta^{\mu \nu} \partial^{2}+\partial^{\mu} \partial^{\nu}\right) A_{r v} \tag{5.182}
\end{equation*}
$$

gives rise to

$$
\begin{equation*}
-i\left(\eta^{\mu v} q^{2}-q^{\mu} q^{v}\right) \delta_{3} \tag{5.183}
\end{equation*}
$$

We compute diagrams with the above rules and stress again that no factors of $Z_{2}$ and $Z_{3}$ appear for external particles because these are already contained in the counterterms.

The above procedure is merely a reorganisation of perturbation theory. The result for an amplitude is the same irrespective of whether the computation is performed

- either starting from $\mathcal{L}$ as given in (5.171) with the original Feynman rules including all $Z$-factors as before,
- or from $\mathcal{L}_{r}^{(1)}+\mathcal{L}_{r}^{(2)}$ in terms of $A_{r}$ and $\psi_{r}$, including counterterm diagrams, and hence without $Z$-factors for external particles.

This reflects the ambiguity in setting up perturbation theory.
As an example, let us compute the relevant 1-loop 1PI diagrams using the renormalized Feynman rules.

- The fully resummed propagator of the renormalised electron field takes the form

$$
\begin{equation*}
\square_{r}=\frac{i}{p-m-\Sigma_{r}(p)} . \tag{5.184}
\end{equation*}
$$

At 1-loop level, i.e. at order $\alpha$ in perturbation theory,

$$
\begin{equation*}
-1 \mathrm{PI}-\left.\right|_{1-\text { loop }}=-\left.i \Sigma_{r}(p)\right|_{1-\text { loop }}=-i \Sigma_{2}^{(1)}(p)+\left.i\left(p \delta_{2}-\delta_{m}\right)\right|_{\alpha} \tag{5.185}
\end{equation*}
$$

Here $-i \Sigma_{2}^{(1)}(p)$ is computed from $\mathcal{L}_{r}^{(1)}$ in terms of $m$ and $e$, i.e. it has the form (5.101) with $e_{0}$ and $m_{0}$ replaced by $e$ and $m$. The term $i\left(p \delta_{2}-\delta_{m}\right)$ is due to the counterterm present in $\mathcal{L}_{r}^{(2)}$. As we will see momentarily, $\delta_{m}$ and $\delta_{2}$ depend on $\alpha$ (as do the remaining counterterms $\delta_{1}$ and $\delta_{3}$ ). At 1-loop level only terms up to and including order $\alpha$ are to be included.

- The photon propgator reads (in Landau gauge for simplicity)

$$
\begin{equation*}
v \sim \sim_{r} \mu=\frac{-i\left(\eta_{\mu \nu}+\frac{q_{\mu} q_{v}}{q^{2}}\right)}{q^{2}\left(1-\Pi_{r}\left(q^{2}\right)\right)} \tag{5.186}
\end{equation*}
$$

At 1-loop level, $\Pi_{r}\left(q^{2}\right)$ is computed via

$$
\begin{equation*}
\left.\sim \Omega \mathrm{PPr}^{2} \sim_{r}\right|_{1-\mathrm{loop}}=\left.i \Pi_{r}^{\mu \nu}(q)\right|_{1-\mathrm{loop}}=i\left(q^{2} \eta^{\mu \nu}-q^{\mu} q^{\nu}\right)\left(\Pi_{2}^{(1)}\left(q^{2}\right)-\left.\delta_{3}\right|_{\alpha}\right) \tag{5.187}
\end{equation*}
$$

with $\Pi_{2}^{(1)}\left(q^{2}\right)$ computed from $\mathcal{L}_{r}^{(1)}$ and $\left.\left.\delta_{3}\right|_{\alpha}\right)$ the counterterm expanded up to order $\alpha$.

- The full vertex is

$$
\begin{equation*}
(\underbrace{p \underbrace{q}_{q} \sim}_{q})_{r}^{p+q}=-i e \Gamma_{r}^{\mu}(p+q, p) \tag{5.188}
\end{equation*}
$$

Again at 1-loop level it takes the form

$$
\begin{equation*}
-\left.i e \Gamma_{r}^{\mu}(p+q, p)\right|_{1-\mathrm{loop}}=-i e \Gamma^{(1) \mu}(p+q, p)-\left.i e \gamma^{\mu} \delta_{1}\right|_{\alpha} \tag{5.189}
\end{equation*}
$$

Finally the 4 counterterm couplings are fixed by the renormalisation conditions as follows:

- Two conditions arise because - ${ }_{r}$ and $v \sim \sim_{r} \mu$ must not involve $Z$-factors at the physical mass poles. This is because we had defined $\psi_{r}$ and $A_{r}$ by this condition, i.e. we had declared that in the renormalised Feynman rules no factors of $Z_{2}$ and $Z_{3}$ appear.
- Two more conditions arise by specifying the meaning of $e$ and $m$. This is arbitrary in principle. One possible choice (out of infinitely many) is that $m$ represents the physical mass and $e$ the physical charge measured at $q=0$.

Combining all of these four conditions translates into the following equations:

$$
\begin{equation*}
(-)_{r} \stackrel{!}{=} \frac{i}{p-m}+\text { terms analytic at } m . \tag{5.190}
\end{equation*}
$$

Therefore we find

$$
\begin{equation*}
\left.\Sigma_{r}(p)\right|_{p m} \stackrel{!}{=} 0 \leftrightarrow m \text { is the physical mass } \tag{5.191}
\end{equation*}
$$

and

$$
\begin{equation*}
1-\left.\frac{\mathrm{d}}{\mathrm{~d} p} \Sigma_{r}(p)\right|_{p=m}=1 \leftrightarrow \text { residue is } 1 \tag{5.192}
\end{equation*}
$$

Identifying $e$ as the physical charge at $q^{2}=0$ amounts to requiring that

$$
\begin{equation*}
(\overbrace{q}^{p})_{r}^{p})_{r}^{p+q}=-i e \Gamma_{r}^{\mu}(p+q, p) \underset{q \rightarrow 0}{\stackrel{!}{\rightarrow}}-i e \gamma^{\mu} . \tag{5.193}
\end{equation*}
$$

Therefore

$$
\begin{equation*}
-i e \Gamma_{r}^{\mu}(p, p) \stackrel{!}{=}-i e \gamma^{\mu} . \tag{5.194}
\end{equation*}
$$

Note that (5.192) has already been used. The last condition is

$$
\begin{equation*}
(\sim \sim \sim)_{r} \stackrel{!}{=} \frac{-i \eta_{\mu \nu}}{q^{2}}+\text { terms analytic at } q=0 \tag{5.195}
\end{equation*}
$$

and thus

$$
\begin{equation*}
1-\Pi_{r}(0) \stackrel{!}{=} 1 \tag{5.196}
\end{equation*}
$$

The conditions (5.191), (5.192), (5.194) and (5.196) are the 4 renormalisation conditions. We can solve these perturbatively, order by order in perturbation theory. At 1-loop order we find for

- condition (5.191)

$$
\begin{equation*}
\Sigma_{2}^{(1)}(m)-\left(m \delta_{2}-\delta_{m}\right) \stackrel{!}{=} 0 \tag{5.197}
\end{equation*}
$$

and therefore

$$
\begin{align*}
\delta_{m} & =m \delta_{2}-\Sigma_{2}^{(1)} \\
& =m \delta_{2}-\frac{\alpha}{2 \pi} \int_{0}^{1} \mathrm{~d} x\left(2-\frac{\epsilon}{2}\right) m-\left(1-\frac{\epsilon}{2}\right) \times m\left(\frac{2}{\epsilon}+\log \frac{M^{2}}{\Delta}\right), \tag{5.198}
\end{align*}
$$

- conditon (5.192)

$$
\begin{equation*}
\delta_{2}=\left.\frac{\mathrm{d}}{\mathrm{~d} p} \Sigma_{2}^{(1)}(p)\right|_{m}, \tag{5.199}
\end{equation*}
$$

- condition (5.194)

$$
\begin{equation*}
\delta_{1}=-\delta F_{1}^{(1)}(0) \tag{5.200}
\end{equation*}
$$

- condition (5.196)

$$
\begin{equation*}
\delta_{3}=\Pi_{2}^{(1)}(0)=-\frac{2 \alpha}{\pi} \int_{0}^{1} \mathrm{~d} x x(1-x)\left(\frac{2}{\epsilon}+\log \frac{M^{2}}{m^{2}}\right) . \tag{5.201}
\end{equation*}
$$

Thus indeed the counterms depend on $\alpha$, as anticipated above.
To conclude the 1-loop renormalised Lagrangian takes the form

$$
\begin{align*}
\mathcal{L}= & \left(-\frac{1}{4} F^{2}+\bar{\psi}(i \not \partial-m) \psi-e \bar{\psi} \gamma^{\mu} \psi A_{\mu}\right) \\
& +\left(-\frac{1}{4} \delta_{3} F^{2}+\bar{\psi}\left(i \delta_{2} \not \partial-\delta_{m}\right) \psi-e \delta_{1} \bar{\psi} \gamma^{\mu} \psi A_{\mu}\right), \tag{5.202}
\end{align*}
$$

where we have dropped the subscript $r$ for the renormalised fields. We summarize:

- $e$ and $m$ are finite quantities to be taken from experiment - the physical coupling and mass.
- $\delta_{1}, \delta_{2}, \delta_{3}$ and $\delta_{m}$ are cutoff-dependent functions, $\delta_{i}=\delta_{i}(e, m, \Lambda)$ in PV (or $\delta_{i}=\delta_{i}(\epsilon, e, m)$ in $\operatorname{dimReg}$ ), which are divergent for $\Lambda \rightarrow \infty$ (or $\epsilon \rightarrow 0$ ). Since we have renormalised to 1-loop order, each $\delta_{i}$ so far is of $O(\alpha)$.
- When computing an amplitude, we take into account both the Feynman rules (5.178) and the counterterms. Since we have computed the counterterms only to order $O(\alpha)$ it is only consistent to compute all other diagrams to that same order, where we must take into account that the $\delta_{i}$ are of order $O(\alpha)$. Thus, as in the computation of the 1PI diagrams, a loop diagram is always accompanied by a counterterm diagram, but to order $\alpha$ no counterterm diagrams appear inside a loop (as this would be order $\alpha^{2}$ ).
- The Feynman rules (5.178) give rise to a divergent expression in terms of $m, e, \Lambda$. The counter term coefficients $\delta_{i}(m, e, \Lambda)$ will precisely cancel these divergences. All physical amplitudes are therefore independent of $\Lambda$ in the end, so we can take $\Lambda \rightarrow \infty$.
- If we wish to compute a quantity to order $\alpha^{2}$ we first need to compute the divergent 1PI terms to order $\alpha^{2}$, and impose again the renormalisation conditions. In this computation we must take into account, in addition to usual 2-loop diagrams from (5.178), also 1-loop diagrams on top of $\mathrm{O}(\alpha)$ counterterms, as well as diagrams involving only counterterms with coefficients expanded to order $\alpha^{2}$. We will discuss this more systematically in a 2-loop example in QFT II.
- The perturbative cancellation of divergences works only because the counterterms are of the same form as the terms in the bare Lagrangian, i.e. they do not introduce any qualitatively new interactions. This is guaranteed because at each order in perturbation theory, no qualitatively new UV divergent 1PI diagrams arise, but the same 1PI diagrams merely receive higher-loop contributions. This leads to a readjustment of the counterterm coefficients order by order such as to absorb the new divergences.


## The renormalisation scheme

- The renormalisation conditions (5.191) and (5.194) are arbitrary: We could define $m$ and $e$ to be any function of the physical mass $m$ or $e$, thereby changing the $\delta_{i}$ accordingly. This ambiguity will cancel in all final amplitudes.
- Conditions (5.192) and (5.196) are fixed by the Feynman rules, which do not contain any Zfactors. However, it is possible to change (5.192) and (5.196) if at the same time we modify the Feynman rules accordingly. I.e. we can include arbitrary $\tilde{Z}$-factors for external states in the Feynman rules, provided we change (5.192) and (5.196). This reflects our freedom in normalising the fields in a manner consistent with the Feynman rules.

The concrete choice of renormalisation condition is called renormalisation scheme.

### 5.10 Infrared divergences

Loop diagrams with massless particles may exhibit infrared divergences from integration over loop momenta $k \rightarrow 0$.

- In QED at 1-loop order the IR divergent 1PI diagrams are


Indeed, we had introduced a small fictitious photon mass $\mu$ to regulate the IR divergences.

- These IR divergences cancel in all cross-sections against another source of IR divergences from radiation of soft photons (Bremsstrahlung). Consider e.g.


As $k \rightarrow 0$ these processes are divergent. This is the case already in classical electrodynamics and called infrared catastrophe of electrodynamics. The reason why this divergence is not that catastrophic after all is that no detector can measure a photon below a certain threshold. Thus we can include all possible Bremsstrahlung photons for $k \rightarrow 0$ to a given process (because as $k \rightarrow 0$ the soft photons cannot be measured. It now so happens that the resulting infrared divergences precisely cancel the IR divergences from the above loop diagrams oder by order in $\alpha$, in the cross-section $\sigma$. For details we refer to Peskin-Schröder 6.1, 6.4, 6.5.

## Chapter 6

## Classical non-abelian gauge theory

### 6.1 Geometric perspective on abelian gauge theory

We had approached $U(1)$ gauge symmetry from the perspective of massless vector fields:

- A consistent Lorentz invariant quantum theory of free massless spin-1 field $A_{\mu}(x)$ must be a gauge theory - see our discussion around equ. (4.67) based on Weinberg I, 8.1.
- At the level of interactions consistency requires that $A_{\mu}(x)$ couples to a conserved current $j^{\mu}(x)$. This is equivalent to the gauging of a global symmetry in the matter sector - see our discussion around (5.54) .

An alternative perspective on $U(1)$ gauge symmetry is as follows:

- Start from a matter theory with a global $U(1)$ symmetry, e.g.

$$
\begin{equation*}
\mathcal{L}=\bar{\psi}(i \not \partial-m) \psi, \tag{6.1}
\end{equation*}
$$

which is invariant under $\psi(x) \mapsto e^{-i e \alpha} \psi(x)$ for constant $\alpha \in \mathbb{R}$.

- In a local QFT it is natural to consider local symmetries, i.e. to promote this to a local transformation

$$
\begin{equation*}
\psi(x) \mapsto e^{-i e \alpha(x)} \psi(x)=: U(x) \psi(x) . \tag{6.2}
\end{equation*}
$$

The logic behind such a modification is that a symmetry transformation at spacetime point $y$ far away from $x$ should not affect the field at $x$. Note that in QFT global symmetries, though, from this perspective, unnatural, are of course fully consistent. By contrast, it is conjectured that in presence of gravity all symmetries must be local and that no global symmetries exist. In any case we are motivated to consider the consequences of the transformation (6.2).

- There is an immediate problem: The ordinary derivative $\partial_{\mu} \psi(x)$, defined via

$$
\begin{equation*}
n^{\mu} \partial_{\mu} \psi(x):=\lim _{\epsilon \rightarrow 0} \frac{1}{\epsilon}\left[\psi\left(x^{\mu}+n^{\mu} \epsilon\right)-\psi(x)\right] \tag{6.3}
\end{equation*}
$$

is not a good object with respect to (6.2) because in (6.3) two objects with very different transformation behaviour under (6.2) appear, i.e.

$$
\begin{equation*}
\psi(x) \mapsto U(x) \psi(x), \quad \text { but } \quad \psi(x+n \epsilon) \mapsto U(x+n \epsilon) \psi(x+n \epsilon) . \tag{6.4}
\end{equation*}
$$

- To define a better notion of derivative we introduce the object $C(x, y)$ - the so-called comparator or Wilson line - such that under (6.2)

$$
\begin{equation*}
C(y, x) \psi(x) \mapsto U(y) C(y, x) \psi(x) . \tag{6.5}
\end{equation*}
$$

Then we can define the covariant derivative $D_{\mu} \psi$ via

$$
\begin{equation*}
n^{\mu} D_{\mu} \psi(x):=\lim _{\epsilon \rightarrow 0} \frac{1}{\epsilon}[\psi(x+n \epsilon)-C(x+n \epsilon, x) \psi(x)] . \tag{6.6}
\end{equation*}
$$

Its transformation under (6.2) is

$$
\begin{equation*}
D_{\mu} \psi(x) \mapsto U(x) D_{\mu} \psi(x) . \tag{6.7}
\end{equation*}
$$

- Let us now construct the Wilson line starting from the requirement (6.5), which implies that under (6.2)

$$
\begin{equation*}
C(y, x) \mapsto U(y) C(y, x) U^{-1}(x) . \tag{6.8}
\end{equation*}
$$

Furthermore we impose $C(y, y)=\mathbb{1}$ for obvious reasons. Note that for the $U(1)$ symmetry under consideration $U(x)$ is a pure phase, and it thus suffices to take $C(y, x)$ as a pure phase. We will soon generalize this.

- Taylor expansion of $C(y, x)$ yields

$$
\begin{equation*}
C(x+\epsilon n, x)=1-i e A_{\mu}(x) \epsilon n^{\mu}+O\left(\epsilon^{2}\right) \tag{6.9}
\end{equation*}
$$

for some vector field $A_{\mu}(x)$. Therefore

$$
\begin{equation*}
D_{\mu} \psi=\partial_{\mu} \psi+i e A_{\mu}(x) \psi(x) \tag{6.10}
\end{equation*}
$$

The transformation behaviour of the Wilson line is

$$
\begin{equation*}
C(x+\epsilon n, x) \mapsto U(x+\epsilon n) C(x+\epsilon n, x) U^{-1}(x) \tag{6.11}
\end{equation*}
$$

and to order $\epsilon$ therefore

$$
\begin{equation*}
1-i e A_{\mu}(x) n^{\mu} \epsilon \mapsto\left(U(x)+\epsilon n^{\mu} \partial_{\mu} U(x)\right)\left(1-i e A_{\mu}(x) n^{\mu} \epsilon\right) U^{-1}(x) . \tag{6.12}
\end{equation*}
$$

Thus the vector field $A_{\mu}$ transforms as

$$
\begin{equation*}
A_{\mu}(x) \mapsto U(x) A_{\mu}(x) U^{-1}(x)+\frac{i}{e} \partial_{\mu} U(x) U^{-1}(x) . \tag{6.13}
\end{equation*}
$$

For $U(x)=e^{-i e \alpha(x)}$ we recover $A_{\mu}(x) \mapsto A_{\mu}(x)+\partial_{\mu} \alpha(x)$ as expected.

- The vector field $A_{\mu}(x)$ is therefore a direct consequence of the existence of a local symmetry. It is called a connection. $A_{\mu}(x)$ is a local field and thus has dynamics in its own right.
- Consider $\left[D_{\mu}, D_{\nu}\right]$ interpreted as acting on $\psi(x)$. Since

$$
\begin{equation*}
D_{\mu} \psi(x) \mapsto U(x) D_{\mu} \psi(x) \tag{6.14}
\end{equation*}
$$

we have

$$
\begin{equation*}
\left[D_{\mu}, D_{\nu}\right] \psi(x) \mapsto U(x)\left[D_{\mu}, D_{\nu}\right] \psi(x)=U(x)\left[D_{\mu}, D_{\nu}\right] U^{-1}(x) U(x) \psi(x) \tag{6.15}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\left[D_{\mu}, D_{\nu}\right] \mapsto U(x)\left[D_{\mu}, D_{\nu}\right] U^{-1}(x) \tag{6.16}
\end{equation*}
$$

- The term

$$
\begin{align*}
{\left[D_{\mu}, D_{\nu}\right] \psi(x) } & =\left[\partial_{\mu}+i e A_{\mu}(x), \partial_{\nu}+i e A_{\nu}(x)\right] \psi(x) \\
& =i e\left(\partial_{\mu} A_{\nu}(x)-\partial_{\nu} A_{\mu}(x)+i e\left[A_{\mu}(x), A_{\nu}(x)\right]\right) \psi(x) \tag{6.17}
\end{align*}
$$

leads to the definition of the field strength or curvature

$$
\begin{equation*}
F_{\mu \nu}:=\frac{1}{i e}\left[D_{\mu}, D_{\nu}\right]=\partial_{\mu} A_{\nu}-\partial_{\nu} A_{\mu}+i e\left[A_{\mu}, A_{\nu}\right] . \tag{6.18}
\end{equation*}
$$

For our $U(1)$ theory with $U(x)=e^{-i e \alpha(x)}$,

$$
\begin{equation*}
F_{\mu \nu}=\partial_{\mu} A_{\nu}-\partial_{\nu} A_{\mu} \tag{6.19}
\end{equation*}
$$

and $F_{\mu \nu}$ is invariant under $U(1)$. Therefore

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu} F^{\mu \nu}+\bar{\psi}\left(i \gamma^{\mu} D_{\mu}-m\right) \psi \tag{6.20}
\end{equation*}
$$

is invariant.

### 6.2 Non-abelian gauge symmetry

We now generalise all of this to a non-abelian symmetry group. Consider a Lie group $H$ of dimension $\operatorname{dim}(H)$. An element $h \in H$ can be written as

$$
\begin{equation*}
H \ni h=\exp \left(-i g \sum_{a=1}^{\operatorname{dim}(H)} \alpha^{a} T^{a}\right), \tag{6.21}
\end{equation*}
$$

where $g \in \mathbb{R}$ takes the role of $e, \alpha^{a} \in \mathbb{R}$ and $T^{a}$ form a basis of the Lie algebra $\operatorname{Lie}(H)$, i.e. the algebra of infinitesimal group transformations. Viewed as an abstract Lie algebra, Lie $(H)$ is determined by the commutation relations

$$
\begin{equation*}
\left[T^{a}, T^{b}\right]=i f^{a b c} T^{c} \tag{6.22}
\end{equation*}
$$

where a sum over $c$ is understood. Without proof we state here that the structure constants $f^{a b c}$ are totally antisymmetric in the indices $a, b, c$ and thus in particular invariant under cyclic permutations of $a, b, c$.
For example:

- The Lie group $H=U(1)$ has $\operatorname{dim}(H)=1$ and its generator is simply $T^{a} \equiv T \in \mathbb{R}$. Therefore $[T, T]=0$ and $H$ is an abelian Lie algebra.
- The Lie group $H=S U(N)$, viewed as an abstract group, is defined as the group of volumeelement preserving linear transformations on $\mathbb{C}^{N}$ which leave the sesqui-linear form

$$
\begin{equation*}
\mathbb{C}^{N} \ni u, v \rightarrow \mathbb{C}:(u, v) \mapsto \sum_{i=1}^{N} \bar{u}_{i} v_{i} \tag{6.23}
\end{equation*}
$$

invariant. We can identify $H$ with the group of $V \in \mathbb{C}^{N, N}$ such that

$$
\begin{equation*}
V^{\dagger}=V^{-1} \quad \text { and } \quad \operatorname{det} V=1 \tag{6.24}
\end{equation*}
$$

by assigning $\forall h \in H$ a matrix $V(h)$ as above. $\operatorname{Lie}(H)$ is then the algebra of $T^{a} \in \mathbb{C}^{N, N}$ such that

$$
\begin{equation*}
T^{a \dagger}=T^{a} \quad \text { and } \quad \operatorname{tr} T^{a}=0 \tag{6.25}
\end{equation*}
$$

The dimension of $\operatorname{Lie}(H)$ is $N^{2}-1$. For instance, the generators of $H=S U(2)$ are typically normalised to be

$$
\begin{equation*}
T^{a}=\frac{1}{2} \sigma^{a} \tag{6.26}
\end{equation*}
$$

with $a=1,2,3$ and $\sigma^{a}$ the Pauli matrices. The structure constants of $S U(2)$ are then

$$
\begin{equation*}
f^{a b c}=\epsilon^{a b c} . \tag{6.27}
\end{equation*}
$$

- Other Lie groups of relevance in physics include $O(N), S O(N), S p(2 N), E_{6}, E_{7}, E_{8}$. For example, $O(N)\left(S O(N)\right.$ ) is the group of (volume element preserving) linear transformations on $\mathbb{R}^{N}$ which leave the bilinear form $u^{i} \delta_{i j} \nu^{j}, i, j=1, \ldots N$, invariant and can be identified with real orthogonal $N \times N$-matrices (of determinant one). $S p(2 N)$ is the group of linear transformations on $\mathbb{R}^{2 N}$ that leave the anti-symmetric symplectic form $u^{i} \omega_{i j} v^{j}, i, j=1, \ldots 2 N$ invariant.

Consider now a matter Lagrangian with matter fields $\psi(x)$ transforming in a unitary representation of $H$ such that $\mathcal{L}$ is invariant under global transformations

$$
\begin{equation*}
\psi(x) \mapsto R(h) \cdot \psi(x), \quad R(h)^{\dagger}=R(h)^{-1} \tag{6.28}
\end{equation*}
$$

with $h \in H$. Two examples are the following:

- Take $H=S U(N)$ and $\psi(x)$ a Dirac spinor field in the fundamental representation, i.e. we consider a $\mathbb{C}^{N}$-valued spinor field such that, suppressing spinor indices,

$$
\forall x: \psi(x) \equiv \psi_{i}(x)=\left(\begin{array}{c}
\psi_{1}(x)  \tag{6.29}\\
\ldots \\
\psi_{N}(x)
\end{array}\right), \quad \quad R(h) \cdot \psi(x) \equiv V_{i j}(h) \psi_{j}(x) .
$$

The Lagrangian

$$
\begin{equation*}
\mathcal{L}=\bar{\psi}(i \not \partial-m) \psi=\sum_{i=1}^{N} \bar{\psi}_{i}(i \not \partial-m) \psi_{i} \tag{6.30}
\end{equation*}
$$

is invariant under a global $S U(N)$ transformation (6.28) because

$$
\begin{equation*}
\bar{\psi}(i \not \partial-m) \psi \mapsto \bar{\psi} R^{\dagger}(h)(i \not \partial-m) R(h) \psi=\bar{\psi} \underbrace{R^{\dagger}(h) R(h)}_{=\mathbb{1}}(i \not \partial-m) \psi . \tag{6.31}
\end{equation*}
$$

By the dimension of a representation we mean the dimension of the vector space in which the matter field takes its value. The (complex) dimension of the fundamental representation of $S U(N)$ is thus $N$.

- Take $H=S U(N)$, but $\psi(x)$ now in the adjoint representation, i.e. we consider now a spinor field valued in the Lie algebra $\operatorname{Lie}(\mathrm{H})$ viewed as a vector space. This means that, with spinor indices suppressed,

$$
\begin{equation*}
\forall x: \psi(x) \equiv \psi_{i j}(x) \in\left\{\mathbb{C}^{N, N} \mid \psi^{\dagger}(x)=\psi(x), \operatorname{tr} \psi(x)=0\right\} \tag{6.32}
\end{equation*}
$$

and the transformation behavior is given by the adjoint action of the Lie group $H$ on its Lie algebra,

$$
\begin{align*}
R(h) \cdot \psi & :=V(h) \psi V(h)^{-1} \equiv V(h)_{i j} \psi_{j k}\left(V^{-1}(k)\right)_{k l}, \\
R^{\dagger}(h) \cdot \psi & :=V^{\dagger}(h) \psi\left(V^{-1}(h)\right)^{\dagger}=V^{-1}(h) \psi V(h) . \tag{6.33}
\end{align*}
$$

The real dimension of the adjoint representation coincides with $\operatorname{dim}(H)$.
Now consider

$$
\begin{equation*}
\mathcal{L}=\operatorname{tr} \bar{\psi}(i \not \partial-m) \psi \equiv \operatorname{tr}\left[\psi_{i j}(i \not \partial-m) \psi_{j l}\right] \equiv \psi_{i j}(i \not \partial-m) \psi_{j i} . \tag{6.34}
\end{equation*}
$$

This is invariant because

$$
\begin{align*}
\operatorname{tr} \bar{\psi}(i \not \partial-m) \psi & \mapsto \operatorname{tr}\left(\left(V \psi V^{-1}\right)^{\dagger} \gamma^{0}(i \not \partial-m) V \psi V^{-1}\right) \\
& =\operatorname{tr}\left(V \bar{\psi} V^{\dagger}(i \not \partial-m) V \psi V^{-1}\right) \\
& =\operatorname{tr}(\bar{\psi} \underbrace{V^{\dagger} V}_{=\mathbb{1}}(i \not \partial-m) \psi \underbrace{V^{-1} V}_{=\mathbb{1}}), \tag{6.35}
\end{align*}
$$

where cyclicity of the trace was used to go from the second to the third line.

We can now repeat all the steps involved in the gauging of $U(1)$ in this more general setting. For definiteness we work with a Dirac spinor field in the fundamental representation of $S U(N)$. Consider the gauge transformation

$$
\begin{equation*}
\psi(x) \mapsto U(x) \psi(x) \tag{6.36}
\end{equation*}
$$

with

$$
\begin{equation*}
U(x)=\exp \left(-i g \sum_{a=1}^{\operatorname{dim}(H)} \alpha^{a}(x) T^{a}\right) \equiv V(h(x)) . \tag{6.37}
\end{equation*}
$$

Now it turns out that the compensator must take values in the representation of $H$. It can therefore be expanded as

$$
\begin{equation*}
C(x+\epsilon n, x)=\mathbb{1}-i g \sum_{a=1}^{\operatorname{dim}(H)} A_{\mu}^{a}(x) T^{a} \epsilon n^{\mu}+O\left(\epsilon^{2}\right) . \tag{6.38}
\end{equation*}
$$

This defines $\operatorname{dim} H$ vector fields $A_{\mu}^{a}(x)$. The object

$$
\begin{equation*}
A_{\mu}(x) \equiv \sum_{a} A_{\mu}^{a}(x) T^{a} \tag{6.39}
\end{equation*}
$$

is then an $N \times N$ matrix-valued vector field. The covariant derivative takes the form

$$
\begin{equation*}
D_{\mu} \psi(x)=\partial_{\mu}(x)+i g \sum_{a} A_{\mu}^{a}(x) T^{a} \psi(x), \tag{6.40}
\end{equation*}
$$

where $T^{a} \psi(x) \equiv T_{i j}^{a} \psi_{j}(x)$. The gauge transformation on $A_{\mu}(x)$ is still

$$
\begin{equation*}
A_{\mu}(x) \mapsto U(x) A_{\mu}(x) U^{-1}(x)+\frac{i}{g} \partial_{\mu} U(x) U^{-1}(x) \tag{6.41}
\end{equation*}
$$

but now $U(x) A_{\mu}(x) U^{-1}(x) \neq A_{\mu}(x)$. Expanding

$$
\begin{equation*}
U(x)=\mathbb{1}-i g \sum_{a} \alpha^{a}(x) T^{a}+O\left(\alpha^{a}(x)^{2}\right) \tag{6.42}
\end{equation*}
$$

one can read off that

$$
\begin{equation*}
A_{\mu}(x) \mapsto A_{\mu}(x)+\partial_{\mu} \alpha^{a}(x) T^{a}-i g \sum_{a} \alpha^{a}(x)\left[T^{a}, A_{\mu}(x)\right] \tag{6.43}
\end{equation*}
$$

where $\left[T^{a}, A_{\mu}(x)\right]=\left[T^{a}, A_{\mu}^{b}(x) T^{b}\right]=A_{\mu}^{b}(x)\left[T^{a}, T^{b}\right]=i f^{a b c} A_{\mu}^{b}(x) T^{c}$. Therefore

$$
\begin{equation*}
A_{\mu}^{c}(x) \mapsto A_{\mu}^{c}(x)+\partial_{\mu} \alpha^{c}(x)+g f^{a b c} \alpha^{a}(x) A_{\mu}^{b}(x) \tag{6.44}
\end{equation*}
$$

The field strength $F_{\mu \nu}=\frac{1}{i g}\left[D_{\mu}, D_{v}\right] \equiv F_{\mu \nu}(x)^{a} T^{a}$ is

$$
\begin{align*}
& F_{\mu \nu}(x)=\partial_{\mu} A_{v}(x)-\partial_{v} A_{\mu}(x)+i g\left[A_{\mu}(x), A_{\nu}(x)\right],  \tag{6.45}\\
& F_{\mu \nu}^{a}(x)=\partial_{\mu} A_{v}^{a}(x)-\partial_{v} A_{\mu}^{a}(x)-g f^{a b c}\left[A_{\mu}^{b}(x), A_{v}^{c}(x)\right] .
\end{align*}
$$

It transforms under a gauge transformation as

$$
\begin{equation*}
F_{\mu \nu}(x) \mapsto U(x) F_{\mu \nu}(x) U^{-1}(x) \tag{6.46}
\end{equation*}
$$

so $F_{\mu \nu}(x)$ is not invariant. Rather $F_{\mu \nu}(x)$ transforms in the adjoint representation. Thus

$$
\begin{equation*}
\operatorname{tr}\left(F_{\mu \nu} F^{\mu \nu}\right) \mapsto \operatorname{tr}(U F_{\mu \nu} \underbrace{U^{-1} U}_{=\mathbb{1}} F^{\mu \nu} U^{-1})=\operatorname{tr}\left(F_{\mu \nu} F^{\mu \nu} U^{-1} U\right)=\operatorname{tr}\left(F_{\mu \nu} F^{\mu \nu}\right) . \tag{6.47}
\end{equation*}
$$

Typically one normalises the generators $T^{a}$ such that

$$
\begin{equation*}
\operatorname{tr} T^{a} T^{b}=\frac{1}{2} \delta^{a b} . \tag{6.48}
\end{equation*}
$$

Then

$$
\begin{align*}
\mathcal{L} & =-\frac{1}{2} \operatorname{tr}\left(F_{\mu \nu} F^{\mu \nu}\right)+\bar{\psi}\left(i \gamma^{\mu} D_{\mu}-m\right) \psi  \tag{6.49}\\
& \equiv-\frac{1}{4} \sum_{a} F_{\mu \nu}^{a} F^{\mu v a}+\bar{\psi}_{i}\left(i \gamma^{\mu} \partial_{\mu}-m\right) \psi_{i}-g \bar{\psi}_{i} \gamma^{\mu} A_{\mu}^{a} T_{i j}^{a} \psi_{j}
\end{align*}
$$

defines the so-called Yang-Mills Lagrangian. The crucial difference to $U(1)$ gauge theory is that the Yang-Mills gauge field exhibits cubic and quartic self-interactions which are contained in the term

$$
\begin{equation*}
-\frac{1}{4} \sum_{a} F_{\mu \nu}^{a} F^{\mu \nu a} \tag{6.50}
\end{equation*}
$$

Diagrammatically, these interactions between the gauge bosons are of the form


In QFT2 we will learn how to quantise such an intrinsically self-interacting theory.

### 6.3 The Standard Model

As a quick application let us briefly sketch that structure of the the Standard Model (SM) of Particle Physics, which is formulated as a Yang-Mills theory specified by the following data:

- The gauge group is

$$
\begin{equation*}
G=S U(3) \times S U(2) \times U(1)_{Y}, \tag{6.51}
\end{equation*}
$$

where $S U(3)$ represents QCD and its 8 gluons, $S U(2)$ describes the weak interactions via $W^{+}$, $W^{-}$and $Z$ and $U(1)_{Y}$ denotes hypercharge. $S U(2)$ and $U(1)_{Y}$ are broken spontaneously to $U(1)_{\text {e.m }}$, whose gauge field is the photon $\gamma$.

- The representations of fermionic matter are given as follows: The SM is a chiral theory, i.e. left- and righthanded fermion fields $\psi_{L} \equiv P_{L} \psi$ and $\psi_{R} \equiv P_{R} \psi$ transform in different representations. The Standard Model comprises 3 families of

|  | $S U(3)$ | $S U(2)$ | $U(1)_{Y}$ |
| ---: | :---: | :---: | :---: |
| lefthanded $Q_{L}$ | $\square$ | $\square \equiv\binom{u}{d}$ | $\frac{1}{6}$ |
| $L$ | $\cdot$ | $\square \equiv\binom{v}{l}$ | $-\frac{1}{2}$ |
| righthanded $u_{R} / d_{R}$ | $\square$ | $\cdot$ | $\frac{2}{3} /-\frac{1}{3}$ |
| $v_{R} / l_{R}$ | $\cdot$ | $\cdot$ | $0 /-1$ |

- In addition there is 1 complex scalar field $\phi(x)$ in representation

$$
\begin{array}{ccc}
S U(3) & S U(2) & U(1)_{Y} \\
\cdot, & \square, & \frac{1}{2}
\end{array}
$$

The interactions are - apart from the Yang-Mills interactions that follow from the above representations - given by

- the Higgs potential

$$
\begin{equation*}
V(\phi)=\mu^{2} \phi^{\dagger} \phi-\lambda\left(\phi^{\dagger} \phi\right)^{2}, \tag{6.52}
\end{equation*}
$$

with $\langle\phi\rangle \neq 0$, which leads to spontaneous breaking of $S U(2) \times U(1)_{Y} \rightarrow U(1)_{\text {e.m. }}$,

- and the gauge invariant Yukawa interactions, very roughly of the form $\phi \bar{\psi}_{i} \psi_{i}$, which yield fermion masses if $\langle\phi\rangle \neq 0$.


## Chapter 7

## Path integral quantisation

### 7.1 Path integral in Quantum Mechanics

The path integral or functional integral formalism ${ }^{1}$ provides a formulation of Quantum Theory completely equivalent to the canonical quantization method exploited to define Quantum Field Theory in the first part of our course. Before turning to the application of this quantization scheme to field theories, we consider path integral methods in non-relativistic Quantum Mechanics. We will start with a quantum mechanical system in the canonical formulation and derive from this the path integral.

### 7.1.1 Transition amplitudes

Consider therefore a 1-particle system with a classical Hamilton function $H(p, q)$. We denote the corresponding quantum mechanical Hamilton operator by $\hat{H}(\hat{p}, \hat{q})$. The Heisenberg picture position operator $\hat{q}_{H}(t)$ has eigenvectors $|q, t\rangle$ with

$$
\begin{equation*}
\hat{q}_{H}(t)|q, t\rangle=q|q, t\rangle \tag{7.1}
\end{equation*}
$$

In the Schrödinger picture we write

$$
\begin{equation*}
\hat{q}_{S}|q\rangle=q|q\rangle \tag{7.2}
\end{equation*}
$$

The Heisenberg and Schrödinger picture eigenstates are related via

$$
\begin{equation*}
|q, t\rangle=e^{i \hat{H}(\hat{p}, \hat{q}) t}|q\rangle \tag{7.3}
\end{equation*}
$$

In Quantum Mechanics we generally consider the transition amplitude

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}} \mid q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\left\langle q_{\mathrm{F}}\right| e^{-i \hat{H}\left(t_{\mathrm{F}}-t_{\mathrm{I}}\right)}\left|q_{\mathrm{I}}\right\rangle \tag{7.4}
\end{equation*}
$$

for which we now derive a path integral representation as follows: The idea is to partition the transition time $T:=t_{\mathrm{F}}-t_{\mathrm{I}}$ into $N+1$ intervals of length

$$
\begin{equation*}
\delta t=\frac{T}{N+1} \tag{7.5}
\end{equation*}
$$

[^31]and to factorise
\[

$$
\begin{equation*}
e^{-i \hat{H} T}=e^{-i \hat{H} \delta t} \ldots e^{-i \hat{H} \delta t} \tag{7.6}
\end{equation*}
$$

\]

Additionally we insert the unit-operator

$$
\begin{equation*}
\mathbb{1}=\int \mathrm{d} q_{k}\left|q_{k}\right\rangle\left\langle q_{k}\right| \tag{7.7}
\end{equation*}
$$

between the factors of $e^{-i \hat{H} \delta t}$ as follows:

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}} \mid q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\lim _{N \rightarrow \infty} \int \prod_{k=1}^{N} \mathrm{~d} q_{k}\left\langle q_{\mathrm{F}}\right| e^{-i \hat{H} \delta t}\left|q_{N}\right\rangle\left\langle q_{N}\right| e^{-i \hat{H} \delta t}\left|q_{N-1}\right\rangle \ldots\left\langle q_{1}\right| e^{-i \hat{H} \delta t}\left|q_{\mathrm{I}}\right\rangle . \tag{7.8}
\end{equation*}
$$

Suppose first that $H(p, q)=f(p)+V(q)$, e.g. with $f(p)=\frac{p^{2}}{2 m}$, and use the Baker-CampbellHausdorff formula

$$
\begin{equation*}
e^{A+B}=e^{A} e^{B} e^{-\frac{1}{2}[A, B]+\cdots} . \tag{7.9}
\end{equation*}
$$

Then we have

$$
\begin{equation*}
e^{-i \hat{H} \delta t}=e^{-i f(\hat{p}) \delta t} e^{-i V(\hat{q}) \delta t} e^{\frac{1}{2}[f(\hat{p}), V(\hat{q})] \delta t^{2}+O\left(\delta t^{3}\right)} . \tag{7.10}
\end{equation*}
$$

The factor quadratic in $\delta t$ is negligible in the limit of $N \rightarrow \infty$, since then $\delta t \rightarrow 0$. Therefore

$$
\begin{align*}
\left\langle q_{k+1}\right| e^{-i \hat{H} \delta t}\left|q_{k}\right\rangle & =\int \mathrm{d} p_{k}\left\langle q_{k+1}\right| e^{-i f(\hat{p}) \delta t}\left|p_{k}\right\rangle\left\langle p_{k}\right| e^{-i V(\hat{q}) \delta t}\left|q_{k}\right\rangle  \tag{7.11}\\
& =\int \mathrm{d} p_{k} e^{-i f\left(p_{k}\right) \delta t} e^{-i V\left(q_{k}\right) \delta t}\left\langle q_{k+1} \mid p_{k}\right\rangle\left\langle p_{k} \mid q_{k}\right\rangle .
\end{align*}
$$

Recalling from Quantum Mechanics that

$$
\begin{equation*}
\left\langle q_{k+1} \mid p_{k}\right\rangle=\frac{1}{\sqrt{2 \pi}} e^{i p_{k} q_{k+1}} \text { and }\left\langle p_{k} \mid q_{k}\right\rangle=\frac{1}{\sqrt{2 \pi}} e^{-i p_{k} q_{k}} \tag{7.12}
\end{equation*}
$$

and taking the limit $\delta t \rightarrow 0$ we obtain for $\left\langle q_{k+1}\right| e^{-i H \delta t}\left|q_{k}\right\rangle$ the expression

$$
\begin{equation*}
\int \frac{\mathrm{d} p_{k}}{2 \pi} e^{-i\left(f\left(p_{k}\right)+V\left(q_{k}\right)\right) \delta t} e^{i\left(q_{k+1}-q_{k}\right) p_{k}} . \tag{7.13}
\end{equation*}
$$

In the limit of infinitesimal intervals we can write this more symmetrically as

$$
\begin{equation*}
\int \frac{\mathrm{d} p_{k}}{2 \pi} e^{-i H\left(p_{k}, \bar{q}_{k}\right) \delta t} e^{i\left(q_{k+1}-q_{k}\right) p_{k}}, \tag{7.14}
\end{equation*}
$$

with

$$
\begin{equation*}
\bar{q}_{k}=\frac{1}{2}\left(q_{k}+q_{k+1}\right) . \tag{7.15}
\end{equation*}
$$

Thus, altogether the transition amplitude is

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}} \mid q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\lim _{N \rightarrow \infty} \int \frac{\mathrm{~d} p_{0}}{2 \pi} \prod_{k=1}^{N} \frac{\mathrm{~d} p_{k} \mathrm{~d} q_{k}}{2 \pi} e^{i \sum_{k=0}^{N}\left[p_{k} \frac{q_{k+1}-q_{k}}{\delta t}-H\left(p_{k}, \bar{q}_{k}\right]\right] \delta t}, \tag{7.16}
\end{equation*}
$$

where $q_{0} \equiv q_{\mathrm{I}}$ and $q_{N+1} \equiv q_{\mathrm{F}}$. Since

$$
\begin{equation*}
p_{k} \frac{q_{k+1}-q_{k}}{\delta t} \rightarrow p_{k} \dot{q}_{k} \quad \text { for } \quad \delta t \rightarrow 0 \tag{7.17}
\end{equation*}
$$

we end up with

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}} \mid q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\int_{q\left(\mathrm{t}_{\mathrm{I}}\right)=q_{\mathrm{I}}}^{q\left(t_{\mathrm{F}}\right)=q_{\mathrm{F}}} \mathcal{D} q(t) \mathcal{D} p(t) e^{i \int_{\mathrm{I}} \mathrm{t}_{\mathrm{F}} \mathrm{~d} t[p \dot{q}-H(p, q)]} . \tag{7.18}
\end{equation*}
$$

The measure $\int \mathcal{D} q(t) \mathcal{D} p(t)$ is defined via the above discretisation with boundary conditions for $q(t)$ as given, whereas $p(t)$ is free at the endpoints. We will say more about the mathematical aspects of this interpretation at the end of this section.
In the above expression it is implicit that we use the prescription (7.14). Indeed the underlying equation

$$
\begin{equation*}
\left\langle q_{k+1}\right| \hat{H}(\hat{p}, \hat{q})\left|q_{k}\right\rangle=\int \frac{\mathrm{d} p_{k}}{2 \pi} e^{i\left(q_{k+1}-q_{k}\right) p_{k}} H\left(p_{k}, \frac{q_{k+1}+q_{k}}{2}\right) \tag{7.19}
\end{equation*}
$$

is certainly correct in the limit $\delta t \rightarrow 0$ if, as assumed, $H(p, q)=f(p)+V(q)$. More generally, the classical Hamiltonian function may contain cross-terms $p^{n} q^{m}$. If $\hat{H}(\hat{p}, \hat{q})$ contains such mixed terms, (7.19) will in general not hold. Note, however, that in the presence of terms $p^{n} q^{m}$ in the classical Hamilton function the quantum mechanical Hamiltonian is not uniquely determined in the first place by the classical one due to ordering issues. Rather, fixing a specific ordering of the mixed terms of $\hat{H}(\hat{p}, \hat{q})$ is part of defining what we mean by the quantum theory - and the process of quantization is not always unique. An operator $\hat{H}(\hat{p}, \hat{q})$ which satisfies equation (7.19), i.e.

$$
\begin{equation*}
H(p, q)=\int \mathrm{d} v e^{i p v}\left\langle q-\frac{v}{2}\right| \hat{H}(\hat{p}, \hat{q})\left|q+\frac{v}{2}\right\rangle \tag{7.20}
\end{equation*}
$$

is called Weyl-ordered. Thus the above derivation of the path integral is, by definition, correct for quantum theories with Weyl-ordered Hamiltonian. As we have seen this trivially includes the case $H(p, q)=f(p)+V(q)$. On the other hand, if $\hat{H}(\hat{p}, \hat{q})$ is not Weyl-ordered, then on the right hand side of (7.18) extra terms appear. This is understood in the sequel - we take the right hand side of (7.18) as an abbreviation for the inclusion of these terms if necessary.

## A note on rigor

A much more detailed and very careful justification of the above procedure leading to equ. (7.18) can be found e.g. in Chapter 1 of Brown, Quantum Field Theory.

We can write equation (7.18) as

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}} \mid q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\int_{q\left(\mathrm{t}_{\mathrm{I}}\right)=q_{\mathrm{I}}}^{q\left(t_{\mathrm{F}}\right)=q_{\mathrm{F}}} \mathcal{D} q(t) \mathcal{D} p(t) e^{i \int_{\mathrm{I}} \mathrm{t}_{\mathrm{F}} \mathrm{~d} t L(p, q)} \tag{7.21}
\end{equation*}
$$

with $\int_{t_{\mathrm{I}}}^{t_{\mathrm{F}}} \mathrm{d} t L(p, q)=: S[p, q]$ and $L(p, q) \equiv p \dot{q}-H(p, q)$. Note that here $L(p, q)$ is really a function of $q$ and $p$, both of which we integrate over separately.

For the special case

$$
\begin{equation*}
H(p, q)=\frac{p^{2}}{2 m}+V(q) \tag{7.22}
\end{equation*}
$$

one can evaluate the $\int \mathcal{D} p$ integral explicitly. This is done rigorously in the above discretisation, where we need to evaluate

$$
\begin{equation*}
\int \frac{\mathrm{d} p_{k}}{2 \pi} e^{i\left(p_{k}\left(q_{k+1}-q_{k}\right)-\delta t \frac{p_{k}^{2}}{2 m}\right)} \tag{7.23}
\end{equation*}
$$

We almost recognize this as a Gaussian integral

$$
\begin{equation*}
\int_{-\infty}^{\infty} \mathrm{d} x e^{-\frac{1}{2} a x^{2}+b x}=\sqrt{\frac{2 \pi}{a}} e^{b^{2}} \quad \quad \operatorname{Re}(a)>0 \tag{7.24}
\end{equation*}
$$

except that the prefactor of the quadratic terms does not yet have a positive part. However, as will be discussed in more detail below, it is mathematically justified to perform an analytic continuation of the integrand by rotating $\delta t$ to the lower complex half-plane. This means that $\delta t$ is replaced by $\delta t(1-i \epsilon)$ with $\epsilon$ small and positive. We then evaluate the integral for $\delta t(1-i \epsilon)$ and take $\epsilon=0$ at the very end of all computations. With this in mind, (7.23) becomes

$$
\begin{equation*}
e^{i \frac{m}{2 \delta t}\left(q_{j+1}-q_{j}\right)^{2}} \sqrt{\frac{-i m}{2 \pi \delta t}} \equiv e^{i \frac{m}{2 \delta t}\left(q_{j+1}-q_{j}\right)^{2}} C \tag{7.25}
\end{equation*}
$$

Then $\frac{m}{2} \dot{q}^{2}-V(q)=L(q, \dot{q})$ remains and we obtain the Feynman-Kac-formula

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}} \mid q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\lim _{N \rightarrow \infty} C^{N+1} \prod_{k=1}^{N} \int \mathrm{~d} q_{k} e^{i \int_{\mathrm{t}_{\mathrm{I}}}^{t_{\mathrm{F}}} \mathrm{~d} t L(q, \dot{q})} \equiv \int_{q\left(t_{\mathrm{I})}=q_{\mathrm{I}}\right.}^{q\left(t_{\mathrm{F}}\right)=q_{\mathrm{F}}} \mathcal{D} q(t) e^{i \int_{\mathrm{I}} \mathrm{t} \mathrm{t} t L(q, \dot{q})} \tag{7.26}
\end{equation*}
$$

We can interpret this as follows: Reinstating $\hbar$, the transition amplitude

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}} \mid q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\left.\int \mathcal{D} q(t) e^{\frac{i}{\hbar} S[q(t), \dot{q}(t)]}\right|_{q\left(t_{\mathrm{I}}\right)=q_{\mathrm{I}}} ^{q(\mathrm{~F})=q_{\mathrm{F}}} \tag{7.27}
\end{equation*}
$$

counts all possible continuous paths from $q_{\mathrm{I}}$ to $q_{\mathrm{F}}$ weighted by $\exp \left[\frac{i}{\hbar} S\right]$. This explains the name path integral. In the classical limit $S[q, \dot{q}] \gg \hbar$ and due to the strongly oscillating phase of the integrand, the right-hand side is dominated by paths for which the action becomes stationary,

$$
\begin{equation*}
\frac{\delta S}{\delta q}=0 \tag{7.28}
\end{equation*}
$$

This way one recovers the classical dynamics as the saddle-point approximation to the full quantum amplitude.

## On the well-definedness of the path integral in Quantum Mechanics:

The above definition of the path integral as the limit $N \rightarrow \infty$ in (7.16) is informal in the mathematical sense due to at least two possible loopholes, which we now close:

- For $N$ finite, we may wonder if each of the integrals (7.14) underlying (7.16) is convergent. For a typical Hamiltonian $H(p, q)$ which is bounded from below, convergence of (7.14) follows by analytic continuation to 'Euclidean time'. Indeed we may view the parameter $\delta t$ in (7.14) as along the (w.l.o.g. positive) real axis in the complex-plane. We can then rotate $H(p, q) \delta t$ to the negative imaginary time-axis by a Wick rotation. Note that the dependence on $p$ and $q$ is unaffected. This Wick rotation is possible in Quantum Mechanics because for an analytic function $H(p, q) \delta t$ we do not encounter any poles in the complex time-plane along the way. Wick rotation then effectively replaces

$$
\begin{equation*}
\delta t \rightarrow-i \delta \tau \tag{7.29}
\end{equation*}
$$

and (7.14) becomes

$$
\begin{equation*}
\int \frac{\mathrm{d} p_{k}}{2 \pi} e^{-H\left(p_{k}, \bar{q}_{k}\right) \delta \tau} e^{i\left(q_{k+1}-q_{k}\right) p_{k}} . \tag{7.30}
\end{equation*}
$$

For instance for $H=\frac{p^{2}}{2 m}+V(q)$, which is the form of $H$ that will be of interest for the extension to field theory, this integral is manifestly convergent as long as the potential $V(q)$ is bounded below.

- The limit $N \rightarrow \infty$ and the ensuing interpretation of (7.16) as an integral over all continuous paths from $q_{I}$ to $q_{F}$ an be made mathematically rigorous - at least for $H=\frac{p^{2}}{2 m}+V(q)$ - by directly passing to the Euclidean theory:
- First rotate $\delta t \rightarrow-i \delta \tau$ already in the classical theory.
- Now compute the Euclidean transition amplitude $\left\langle q_{\mathrm{F}}\right| e^{-\hat{H} \tau}\left|q_{\mathrm{I}}\right\rangle$. At least for the physically most interesting theory $H=\frac{p^{2}}{2 m}+V(q)$ one can rigorously prove that a mathematically well-defined measure exists on the space of continuous paths from $q_{I}$ to $q_{F}$ - the so-called Wiener measure. The existence of this measure for $H=\frac{p^{2}}{2 m}+V(q)$ depends on the Gaussian suppression factor in the Euclidean theory and furthermore on the crucial factor $C=\sqrt{\frac{-i m}{2 \pi \delta t}}$ that appears in evaluating the momentum integral. ${ }^{2}$
- In the end we rotate the Euclidean expression back to real time by making the replacement $\tau \rightarrow i t$ in all formulae.

To conclude, the path-integral in Quantum mechanics is mathematically well-defined - at least for $H=\frac{p^{2}}{2 m}+V(q)$, which is the form of the Hamiltonian that we will need for generalizations to field theory. Note that despite the lack of rigour in the derivation of (7.16), for a number of interesting

[^32]theories such as the harmonic oscillator the path integral in this form can be solved and convergent expressions can be obtained in a very explicit manner. This will be demonstrated in the tutorials. Crucial in this context is again a suitable normalization of the position space path integral by carefully taking into account the factor $C=\sqrt{\frac{-i m}{2 \pi \delta t}}$ in the measure in (7.26).

### 7.1.2 Correlation functions

Consider now the matrix element

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}}\right| \hat{q}_{H}\left(t_{1}\right)\left|q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle \tag{7.31}
\end{equation*}
$$

of the Heisenberg operator $q_{H}\left(t_{1}\right)$ for $t_{\mathrm{F}}>t_{1}>t_{\mathrm{I}}$. We can transform this to the Schrödinger picture as

$$
\begin{equation*}
\hat{q}_{H}\left(t_{1}\right)=e^{i \hat{H} t_{1}} \hat{q}_{S} e^{-i \hat{H} t_{1}} \tag{7.32}
\end{equation*}
$$

and obtain

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}\right| e^{-i \hat{H}\left(t_{\mathrm{F}}-t_{1}\right)} \hat{q}_{S} e^{-i \hat{H}\left(t_{1}-t_{\mathrm{I}}\right)}\left|q_{\mathrm{I}}\right\rangle \tag{7.33}
\end{equation*}
$$

We proceed as before in the computation of the transition amplitude by factorising the time evolution operator and inserting the identity multiple times. The operator $\hat{q}_{S}$ now acts on $|q\rangle_{K}\left\langle\left. q\right|_{K}\right.$ corresponding to the time $t_{1}$,

$$
\begin{equation*}
\hat{q}_{S}|q\rangle_{K}\left\langle\left.\left. q\right|_{K}\right|_{t_{1}}=q\left(t_{1}\right) \mid q\right\rangle_{K}\left\langle\left.\left. q\right|_{K}\right|_{t_{1}} .\right. \tag{7.34}
\end{equation*}
$$

The matrix elements thus becomes

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}}\right| \hat{q}_{H}\left(t_{1}\right)\left|q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\int_{q\left(t_{\mathrm{I}}\right)=q_{\mathrm{I}}}^{q\left(t_{\mathrm{F}}\right)=q_{\mathrm{F}}} \mathcal{D} q(t) \mathcal{D} p(t) q\left(t_{1}\right) e^{i S[q, p]} . \tag{7.35}
\end{equation*}
$$

Likewise for $t_{\mathrm{F}}>t_{1}>t_{2}>t_{\mathrm{I}}$ we find

$$
\begin{equation*}
\int \mathcal{D} q \mathcal{D} p q\left(t_{1}\right) q\left(t_{2}\right) e^{i S[q, p]}=\left\langle q_{\mathrm{F}}, t_{\mathrm{F}}\right| \hat{q}_{H}\left(t_{1}\right) \hat{q}_{H}\left(t_{2}\right)\left|q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle, \tag{7.36}
\end{equation*}
$$

while for $t_{\mathrm{F}}>t_{2}>t_{1}>t_{1}$ we would get

$$
\begin{equation*}
\int \mathcal{D} q \mathcal{D} p q\left(t_{1}\right) q\left(t_{2}\right) e^{i S[q, p]}=\left\langle q_{\mathrm{F}}, t_{\mathrm{F}}\right| \hat{q}_{H}\left(t_{2}\right) \hat{q}_{H}\left(t_{1}\right)\left|q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle . \tag{7.37}
\end{equation*}
$$

Thus time-ordering appears automatically:

$$
\begin{equation*}
\left\langle q_{\mathrm{F}}, t_{\mathrm{F}}\right| T \hat{q}_{H}\left(t_{1}\right) \hat{q}_{H}\left(t_{2}\right)\left|q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=\int_{q\left(t_{\mathrm{t}_{1}}\right)=q_{\mathrm{I}}}^{q\left(t_{\mathrm{F}}\right)=q_{\mathrm{F}}} \mathcal{D} q \mathcal{D} p q\left(t_{1}\right) q\left(t_{2}\right) e^{i S[q, p]} . \tag{7.38}
\end{equation*}
$$

We will need above expressions with initial and final states given by $|\Omega\rangle$ - the vacuum of the interacting theory. To compute this expression most simply we use the same trick employed in the context of
the interaction picture derivation of correlation functions in QFT I. Consider an eigenbasis of the Hamiltonian with

$$
\begin{equation*}
\hat{H}|n\rangle=E_{n}|n\rangle \text { and } \hat{H}|\Omega\rangle=E_{\Omega}|\Omega\rangle \tag{7.39}
\end{equation*}
$$

and suppose without loss of generality that $E_{\Omega}=0$ (otherwise we shift $\hat{H}$ by a constant). Then $\forall\left|q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle$ we can write

$$
\begin{equation*}
\left|q_{\mathrm{I}}, t_{\mathrm{I}}\right\rangle=e^{i \hat{H} t_{\mathrm{I}}}\left|q_{\mathrm{I}}\right\rangle=\sum_{n} e^{i \hat{H} t_{\mathrm{I}}}|n\rangle\left\langle n \mid q_{\mathrm{I}}\right\rangle=|\Omega\rangle\left\langle\Omega \mid q_{\mathrm{I}}\right\rangle+\sum_{|n\rangle \neq|\Omega\rangle} e^{i E_{n} t_{\mathrm{I}}}|n\rangle\left\langle n \mid q_{\mathrm{I}}\right\rangle \tag{7.40}
\end{equation*}
$$

Let us now replace $t_{\mathrm{I}} \rightarrow t_{\mathrm{I}}(1-i \epsilon)$ and take $t_{\mathrm{I}} \rightarrow-\infty$, with $\epsilon$ small and positive. In this limit all terms with $|n\rangle \neq|\Omega\rangle$ are exponentially suppressed. Thus

$$
\begin{equation*}
\langle\Omega \mid \Omega\rangle=\left(\left\langle q_{\mathrm{F}} \mid \Omega\right\rangle\left\langle\Omega \mid q_{\mathrm{I}}\right\rangle\right)^{-1} \lim _{T \rightarrow \infty(1-i \epsilon)}\left\langle q_{\mathrm{F}}, T \mid q_{\mathrm{I}},-T\right\rangle \tag{7.41}
\end{equation*}
$$

for all $\left|q_{I}\right\rangle,\left|q_{\mathrm{F}}\right\rangle$ with non-zero overlap with $|\Omega\rangle$. Therefore

$$
\begin{equation*}
\langle\Omega| T \prod_{i} \hat{q}_{H}\left(t_{i}\right)|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{\int \mathcal{D} q(t) \mathcal{D} p(t) e^{i \int^{-T} \mathrm{~d} t L(p, q)} \prod_{i} q\left(t_{i}\right)}{\int \mathcal{D} q(t) \mathcal{D} p(t) e^{i \int^{-T} \mathrm{~d} t L(p, q)}} \tag{7.42}
\end{equation*}
$$

Note that

- the boundaries in $\int \mathcal{D} q(t)$ are arbitrary as long as $\left\langle\Omega \mid q_{\mathrm{I}}\right\rangle \neq 0 \neq\left\langle q_{\mathrm{F}} \mid \Omega\right\rangle$ and
- the denominator ensures that $\langle\Omega| \mathbb{1}|\Omega\rangle=1$.


### 7.2 The path integral for scalar fields

We now apply the path integral approach to a field theory. Consider, to begin with, a classical real scalar field theory with classical Hamiltonian

$$
\begin{equation*}
H=\int \mathrm{d}^{3} x \mathcal{H}(\phi(x), \Pi(x)) \tag{7.43}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{H}(\phi(x), \Pi(x))=\frac{1}{2} \Pi^{2}(x)+\frac{1}{2}(\nabla \phi(x))^{2}+V(\phi(x)) . \tag{7.44}
\end{equation*}
$$

Note that all non-hatted quantities are classical fields. One defines the quantum fields $\hat{\phi}_{H}(t, \vec{x})$ in the Heisenberg picture and its eigenstates via

$$
\begin{equation*}
\hat{\phi}_{H}(t, \vec{x})|\phi(\vec{x}), t\rangle=\phi(t, \vec{x})|\phi(\vec{x}), t\rangle . \tag{7.45}
\end{equation*}
$$

Our task is thus to compute the transition amplitudes between two such states

$$
\begin{equation*}
\left\langle\phi_{\mathrm{F}}(\vec{x}), t_{\mathrm{F}} \mid \phi_{\mathrm{I}}(\vec{x}), t_{\mathrm{I}}\right\rangle \equiv\left\langle\phi_{\mathrm{F}}(\vec{x})\right| e^{i \hat{H}\left(t_{\mathrm{F}}-t_{\mathrm{I}}\right)}\left|\phi_{\mathrm{I}}(\vec{x})\right\rangle . \tag{7.46}
\end{equation*}
$$

One proceeds by a careful regularization procedure as follows:

- Introduce an infrared (IR) cutoff by considering the theory in a finite volume $V$ as well as an ultraviolet (UV) cutoff by discretizing $V$. I.e. we consider the theory on a finite 3 -dimensional lattice $L$ with lattice spacing $a$ and parametrise the lattice sites by the discrete vectors

$$
\begin{equation*}
\vec{x}=a\left(n_{1}, n_{2}, n_{3}\right), \quad n_{i} \in\{-N, \ldots, N\} . \tag{7.47}
\end{equation*}
$$

The Hamiltonian of this discretized theory is

$$
\begin{equation*}
\hat{H}=\sum_{\vec{x} \in L}\left(\frac{1}{2} \hat{\pi}^{2}(\vec{x})+\frac{1}{2} \sum_{\vec{j}}\left(\frac{\hat{\phi}(\vec{x}+\vec{j})-\hat{\phi}(\vec{x})}{a}\right)^{2}+V(\hat{\phi}(\vec{x}))\right), \tag{7.48}
\end{equation*}
$$

where the sum over $\vec{j}$ is a suitable sum over nearest neighbors such as to reproduce the spatial derivative.

- The IR and UV regularized theory thus corresponds to a 3-dimensional Quantum Mechanical system with a finite number of degrees of freedom. Thus we can use the path-integral formula for the computation of the transition amplitude. By Wick rotation this is completely welldefined - including the transition to continuous time (but with space discrete).
- In the end we remove the IR and the UV cutoff in space. It is this continuum limit which introduces both IR and UV divergences. These must be taken care of by carefully implementing renormalisation as we will see in more detail later. In a properly renormalized theory, all physical quantities computed via the path integral (in particular correlation functions) are finite even after taking the continuum limit.

As an abbreviation of this procedure we can write the above transition amplitude as

$$
\begin{align*}
\left\langle\phi_{\mathrm{F}}(\vec{x}), t_{\mathrm{F}} \mid \phi_{\mathrm{I}}(\vec{x}), t_{\mathrm{I}}\right\rangle & \equiv\left\langle\phi_{\mathrm{F}}(\vec{x})\right| e^{i \hat{H}\left(t_{\mathrm{F}}-t_{\mathrm{I}}\right)}\left|\phi_{\mathrm{I}}(\vec{x})\right\rangle \\
& =\int_{\phi\left(\vec{x}, t_{1}\right)=\phi_{\mathrm{I}}(\vec{x})}^{\left.\phi_{\mathrm{t}}, t_{\mathrm{F}}\right)=\phi_{\mathrm{F}}(\vec{x})} \mathcal{D} \phi(x) \mathcal{D} \Pi(x) e^{i \int_{\mathrm{I}_{\mathrm{F}}}^{t_{\mathrm{F}}} \mathrm{~d}^{4} x[\Pi(x) \dot{\phi}(x)-\mathcal{H}(\phi(x), \Pi(x))]} . \tag{7.49}
\end{align*}
$$

This is a formal expression meant in the above sense. Inserting the specific Hamiltonian (7.43) in the exponent shows that the $\Pi$ integral is Gaussian, since $\dot{\phi} \equiv \Pi$. The Gaussian integral can be performed explicitly as in the finite-dimensional case. A rigorous derivation can be given by discretizing the path integral and taking the continuum limit in the end. Altogether we have

$$
\begin{equation*}
\left\langle\phi_{\mathrm{F}}(\vec{x}), t_{\mathrm{F}} \mid \phi_{\mathrm{I}}(\vec{x}), t_{\mathrm{I}}\right\rangle=\int_{\phi\left(\vec{x}, t_{\mathrm{I}}\right)=\phi_{\mathrm{I}}(\vec{x})}^{\phi\left(\vec{x}, t_{\mathrm{F}}\right)=\phi_{\mathrm{F}}(\vec{x})} \mathcal{D} \phi(x) e^{i \int_{\mathrm{I}_{\mathrm{I}}} i \mathrm{~d}^{\mathrm{t}^{4}} x \mathcal{L}(\phi(x))}, \tag{7.50}
\end{equation*}
$$

where the extra factors picked up upon performing the Gaussian integral are, by definition, included in the measure. Since we will normalize the relevant expressions in a suitable manner momentarily
this has no importance.
Using again the trick from (7.41) yields, up to a factor of proportionality,

$$
\begin{equation*}
\langle\Omega \mid \Omega\rangle \sim \lim _{T \rightarrow \infty(1-i \epsilon)} \mathcal{D} \phi e^{i \int^{T} \mathrm{~d}^{4} x \mathcal{L}(\phi(x))} \tag{7.51}
\end{equation*}
$$

with arbitrary boundary conditions with non-zero overlap with $|\Omega\rangle$. This yields the path-integral master formula for the computation of quantum correlation functions

$$
\begin{align*}
G\left(x_{1}, \ldots, x_{n}\right) & \equiv\langle\Omega| T \prod_{i} \hat{\phi}_{H}\left(x_{i}\right)|\Omega\rangle \\
& =\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{\int \mathcal{D} \phi(x) \phi\left(x_{1}\right) \ldots \phi\left(x_{n}\right) e^{i \int_{-T}^{T} \mathrm{~d}^{4} x \mathcal{L}(\phi(x))}}{\int \mathcal{D} \phi(x) e^{i \int_{-T}^{T} \mathrm{~d}^{4} x \mathcal{L}(\phi(x))}}, \tag{7.52}
\end{align*}
$$

where the denominator ensures again correct normalisation. Note in particular that all universal constant prefactors from the Gaussian integral cancel as promised.

## On Euclidean Field Theory

As pointed out, the appearance of possible divergences in the path integral is related to the ultimate removal of the spatial cutoff in the defining expression for the path-integral. In a perturbatively renormalisable theory, the action can be renormalized order by order such that the perturbatively expanded correlation functions are finite in the continuum limit, order by order. Furthermore, as we will see, to set up perturbation theory it is sufficient to perform formal manipulations with the formal expression of the path-integral without worrying about any convergence issues. More ambitiously, one might hope to directly evaluate the path integral in an exact, non-perturbative manner. One way to do this in practice is Lattice Field Theory, where one keeps, however, the spatial cutoff finite and computes the integral numerically.
In this approach, but also in more general contexts, if one is interested in explicit evaluations of the path integral it is, as in Quantum Mechanics, useful to exploit the better convergence properties of the appearing integrals in the Euclidean theory. The idea is to compute the correlation functions as functions on $\mathbb{R}^{4}$ (as opposed to $\mathbb{R}^{1,3}$ ) with coordinates

$$
\begin{equation*}
x_{E}=\left(x^{1}, x^{2}, x^{3}, x^{4}=i x^{0}\right) . \tag{7.53}
\end{equation*}
$$

The Wick rotation $x^{4}=i x^{0}$ implies that we should replace

$$
\begin{equation*}
e^{i \int d^{4} x \mathcal{L}(\phi(x))} \rightarrow e^{-\int d^{4} x_{E} \mathcal{L}\left(\phi\left(x_{E}\right)\right)}=e^{-S_{E}} . \tag{7.54}
\end{equation*}
$$

It is the suppression of the path integral by $e^{-S_{E}}$ which is responsible for its improved convergence properties. The correlation functions in the Euclidean theory are called Schwinger functions, as opposed to the correlation functions in Minkowski space, called Wightman functions. This way, one can
compute a correlation function of a theory in $\mathbb{R}^{1,3}$ by first computing its associated Schwinger function in the Euclidean theory and then Wick rotating back to $\mathbb{R}^{1,3}$ - provided one does not encounter any poles in the complex $x_{E}$-plane as one undoes the Wick rotation.
The Osterwalder-Schrader theorem gives conditions under which a mathematically precise definition of a Lorentzian QFT (in the sense that it satisfies the so-called Wightman axions) via analytic continuation of the Euclidean path integral is possible. In particular, the Schwinger functions must satisfy a property known as reflection positivity. Unfortunately, to date the only interacting field theories for which the Osterwalder-Schrader conditions have been rigorously proven to be satisfied are certain theories in 2 and 3 dimensions. ${ }^{3}$ These theories can be viewed as mathematically well-defined quantum field theories. Establishing such a definition of Quantum Yang-Mills theory in $\mathbb{R}^{1,3}$ is a 1-million-dollar problem.

## On wave functionals

We can conclude that there are two equivalent ways to define a quantum field theory starting from a classical theory:

## - Canonical formalism:

First promote the classical field $\phi(x)$ to a quantum operator $\hat{\phi}_{H}(x)$. Then define the Hilbert space of states and derive the correlators in the interaction picture as given by

$$
\begin{equation*}
\langle\Omega| T \prod_{i} \hat{\phi}_{H}\left(x_{i}\right)|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{\langle 0| T \prod_{i} \hat{\phi}_{I}\left(x_{i}\right) e^{-i} \int_{-T}^{T} \mathrm{~d} t H_{I}(t)}{\langle 0\rangle} \underset{\langle 0| T e^{-i} \int_{-T}^{T} \mathrm{~d} H_{l}(t)}{ }|0\rangle . \tag{7.55}
\end{equation*}
$$

This is the Gell-Mann-Low formula, which we derived in section 2.5 in QFT I.

## - Path integral formalism:

Alternatively we can begin not with the concept of a quantum operator, but with the concept of a physical state as a wave functional of the classical fields. This is a generalization of the well-familiar correspondence between states and wave functions (as functions on the space of classical observables) in Quantum Mechanics. Consider for instance the Schrödinger representation in Quantum Mechanics, where the correspondence between a state localised at $q_{0}$ and its wavefunction is

$$
\begin{equation*}
\left|q_{0}, t\right\rangle \leftrightarrow \psi(q(t))=\delta\left(q(t)-q_{0}(t)\right) . \tag{7.56}
\end{equation*}
$$

Similarly, in Quantum Field Theory we define a state $\left|\phi_{0}(\vec{x}, t)\right\rangle$ via a functional $\Psi(\phi(\vec{x}, t))$ (as a map from the space of classical observables - the classical fields - to $\mathbb{C}$ ) as

$$
\begin{equation*}
\left|\phi_{0}(\vec{x}, t)\right\rangle \leftrightarrow \Psi[\phi]=\delta\left[\phi(\vec{x}, t)-\phi_{0}(\vec{x}, t)\right] . \tag{7.57}
\end{equation*}
$$

[^33]The next step is then to define the concept of a quantum operator $\hat{\phi}_{H}(x)$ by specifying its matrix elements with a basis of states as

$$
\begin{equation*}
\left\langle\phi_{1}(\vec{x}), t_{1}\right| \hat{\phi}_{H}(x)\left|\phi_{2}(\vec{x}), t_{2}\right\rangle=\int_{\phi\left(\vec{x}, t_{1}\right)=\phi_{1}(\vec{x})}^{\phi\left(\vec{x}, t_{2}\right)=\phi_{2}(\vec{x})} \mathcal{D} \phi(x) \mathcal{D} \Pi(x) e^{i \int_{t_{1}}^{t_{1}} \mathrm{~d}^{4} x[\Pi \dot{\phi}-\mathcal{H}]} \tag{7.58}
\end{equation*}
$$

and more generally the correlators $G\left(x_{1}, \ldots, x_{n}\right)$ via path integral of classical fields. In particular the path integral is a convenient way to express the vacuum wave functional as the overlap

$$
\begin{equation*}
\Psi_{0}[\phi]=\langle\phi(\vec{x}), t \mid \Omega\rangle=\lim _{t_{1} \rightarrow-\infty(1-i \epsilon)} \int \mathcal{D} \phi\left(\vec{x}, t^{\prime}\right) \mathcal{D} \Pi\left(\vec{x}, t^{\prime}\right) e^{i \int_{t_{1}}^{t} \mathrm{~d}^{\prime} \mathrm{d}^{3} x[\Pi \dot{\phi}-\mathcal{H}]} \tag{7.59}
\end{equation*}
$$

where $\left.\phi\left(\vec{x}, t^{\prime}\right)\right|_{t^{\prime}=t}=\phi(\vec{x}, t)$, but $\phi\left(\vec{x}, t_{\mathrm{I}}\right)$ is arbitrary (with non-zero overlap with $|\Omega\rangle$ ). ${ }^{4}$

### 7.3 Generating functional for correlation functions

A functional maps a function - e.g. the classical field $\phi(x)$ - to a $\mathbb{C}$-number, e.g.

$$
\begin{equation*}
\phi(x) \mapsto S[\phi(x)]=\int \mathrm{d}^{4} x \mathcal{L}(\phi(x)) \tag{7.60}
\end{equation*}
$$

Consider the functional

$$
\begin{equation*}
Z[J]:=\int \mathcal{D} \phi e^{i S[\phi]+i \int \mathrm{~d}^{4} x J(x) \phi(x)} \tag{7.61}
\end{equation*}
$$

with some source-term $J(x)$. Then we can rewrite this as

$$
\begin{align*}
Z[J]= & \int \mathcal{D} \phi e^{i S[\phi]}\left(1+i \int \mathrm{~d}^{4} x \phi(x) J(x)+\frac{i^{2}}{2!} \int \mathrm{d}^{4} x_{1} \mathrm{~d}^{4} x_{2} \phi\left(x_{1}\right) \phi\left(x_{2}\right) J\left(x_{1}\right) J\left(x_{2}\right)+\ldots\right) \\
= & \int \mathcal{D} \phi e^{i S[\phi]}+i \int \mathrm{~d}^{4} x J(x) \int \mathcal{D} \phi \phi(x) e^{i S[\phi]}  \tag{7.62}\\
& +\frac{i^{2}}{2!} \int \mathrm{d}^{4} x_{1} \mathrm{~d}^{4} x_{2} J\left(x_{1}\right) J\left(x_{2}\right) \int \mathcal{D} \phi \phi\left(x_{1}\right) \phi\left(x_{2}\right) e^{i S[\phi]}+\ldots
\end{align*}
$$

and thus

$$
\begin{equation*}
\frac{Z[J]}{Z[0]}=\sum_{n=0}^{\infty} \frac{i^{n}}{n!} \int \mathrm{d}^{4} x_{1} \ldots \int \mathrm{~d}^{4} x_{n} J\left(x_{1}\right) \ldots J\left(x_{n}\right) G\left(x_{1}, \ldots, x_{n}\right) . \tag{7.63}
\end{equation*}
$$

To solve this formula for $G\left(x_{1}, \ldots, x_{n}\right)$ we need the tools of functional calculus.

### 7.3.1 Functional calculus

We can think of $\phi(x)$ as an element of an $\infty$-dimensional vector space and consider the following formal analogies:

[^34]
## - Finite-dimensional vector calculus

A vector is given as $u \equiv u_{i}$. We can map two real vectors to $\mathbb{R}$ via the inner product

$$
\begin{equation*}
u_{i} v_{i} \equiv u^{T} v \equiv u \cdot v \tag{7.64}
\end{equation*}
$$

Linear operators are defined via

$$
\begin{equation*}
v_{i}=A_{i j} u_{j} \equiv(A \cdot v)_{i} . \tag{7.65}
\end{equation*}
$$

The identity operator is

$$
\begin{equation*}
A_{i j}=\delta_{i j} \tag{7.66}
\end{equation*}
$$

such that the inverse of a linear operator satisfies

$$
\begin{equation*}
A_{i j} A_{j k}^{-1}=\delta_{i k} . \tag{7.67}
\end{equation*}
$$

## - Infinite dimensional vector calculus (=functional calculus)

In this case the vector is the classical field $\phi(x)$, with the dependence on the continuous parameter $x$ formally representing the previous index $i$. The inner product becomes

$$
\begin{equation*}
\int \mathrm{d} x \phi(x) \chi(x) \equiv \phi \cdot \chi \tag{7.68}
\end{equation*}
$$

A linear operator now acts as

$$
\begin{equation*}
\chi(x)=\int \mathrm{d} y K(x, y) \phi(y) \equiv(K \cdot \phi)(x) \tag{7.69}
\end{equation*}
$$

where $K(x, y)$ is called the integral kernel. The identity operator is

$$
\begin{equation*}
K(x, y)=\delta(x-y), \tag{7.70}
\end{equation*}
$$

such that the inverse operator obeys

$$
\begin{equation*}
\int \mathrm{d} y K(x, y) K^{-1}(y, z)=\delta(x-z) \tag{7.71}
\end{equation*}
$$

To set up a differential calculus we first consider a function of a finite-dimensional vector and perform a Taylor expansion

$$
\begin{equation*}
f\left(u_{i}+h_{i}\right)=f\left(u_{i}\right)+\frac{\partial f}{\partial u_{i}} h_{i}+O\left(h^{2}\right) . \tag{7.72}
\end{equation*}
$$

This defines the derivative $\frac{\partial f}{\partial u_{i}}$. Formally we can generalise this definition of a derivative to a functional $F[\varphi(x)]$ by expanding

$$
\begin{equation*}
F[\varphi(x)+h(x)]=F[\varphi(x)]+\int \mathrm{d} x \frac{\delta F}{\delta \varphi(x)} h(x)+O\left(h^{2}\right) . \tag{7.73}
\end{equation*}
$$

The quantity $\frac{\delta F}{\delta \varphi(x)}$ is called functional derivative. It obeys the same rules as the finite-dimensional derivatives. In particular it is useful to consider the analogies

$$
\begin{align*}
\frac{\partial}{\partial u_{i}} u_{j}=\delta_{i j} & \leftrightarrow \frac{\delta}{\delta \varphi(x)} \varphi(y)=\delta(x-y),  \tag{7.74}\\
\frac{\partial}{\partial u_{i}}\left(u_{j} k_{j}\right)=\delta_{i j} k_{j}=k_{i} & \leftrightarrow \frac{\delta}{\delta \varphi(x)} \int \mathrm{d} y \varphi(y) \chi(y) \equiv \frac{\delta}{\delta \varphi(x)}(\varphi \cdot \chi)=\chi(x) . \tag{7.75}
\end{align*}
$$

Further useful rules are

$$
\begin{align*}
& \frac{\delta}{\delta \varphi(x)} \int \mathrm{d} y\left(\partial_{y} \varphi(y)\right) \chi(y)=\frac{\delta}{\delta \varphi(x)}\left[-\int \mathrm{d} y \varphi(y) \partial_{y} \chi(y)\right]=-\partial_{x} \chi(x)  \tag{7.76}\\
& \frac{\delta}{\delta \varphi(x)} e^{i \int \mathrm{~d} y \varphi(y) \chi(y)}=i \chi(x) e^{i \int \mathrm{~d} y \varphi(y) \chi(y)} .
\end{align*}
$$

We will oftentimes leave the actual integrals implicit to save some writing. The last equation then reads

$$
\begin{equation*}
\frac{\delta}{\delta \varphi(x)} e^{i \varphi \cdot \chi}=i \chi(x) e^{i \varphi \cdot \chi} \tag{7.77}
\end{equation*}
$$

Going back to (7.63) it is straightforward to see that

$$
\begin{equation*}
G\left(x_{1}, \ldots, x_{n}\right)=\left.\frac{1}{Z[0]}\left(\frac{\delta}{i \delta J\left(x_{1}\right)}\right) \ldots\left(\frac{\delta}{i \delta J\left(x_{n}\right)}\right) Z[J]\right|_{J=0} \tag{7.78}
\end{equation*}
$$

$Z[J]$ is called generating functional for the Green's functions.

### 7.4 Free scalar field theory

We consider the free scalar field action

$$
\begin{align*}
S_{0}[\phi] & =\int \mathrm{d}^{4} x\left(\frac{1}{2} \partial_{\mu} \phi \partial^{\mu} \phi-\frac{1}{2} m_{0}^{2} \phi^{2}\right) \\
& =\int \mathrm{d}^{4} x\left(-\frac{1}{2} \phi(x)\left(\partial_{x}^{2}+m_{0}^{2}\right) \phi(x)\right) . \tag{7.79}
\end{align*}
$$

In the path integral formula of $Z[J]$ the contour is taken along $t(1-i \epsilon)$ in order to project the initial and final states to the vacuum state. For the free theory one can prove that alternatively one can shift not the time variable into the complex plane, but make the replacement

$$
\begin{equation*}
\mathcal{H} \rightarrow \mathcal{H}-i \epsilon \phi^{2}(x), \quad \text { i.e. } \quad m_{0}^{2} \rightarrow m_{0}^{2}-i \epsilon . \tag{7.80}
\end{equation*}
$$

Since this prescription works for the free theory, it does so also for perturbative interactions as long as the interaction does not deform the theory too far away from the free theory. Non-perturbatively, however, the replacement (7.80) cannot in general be proven to serve our purpose of projecting the
initial and final states to the vacuum and one must stick with the prescription $t \rightarrow t(1-i \epsilon)$.
With (7.80) the partition function for the free theory takes the form

$$
\begin{align*}
Z_{0}[J] & =\int \mathcal{D} \phi e^{\frac{i}{2} \int \mathrm{~d}^{4} x \mathrm{~d}^{4} y \phi(x) K(x, y) \phi(y)+i \int \mathrm{~d}^{4} x \phi(x) J(x)}  \tag{7.81}\\
& =\int \mathcal{D} \phi e^{\frac{i}{2} \phi \cdot(K \phi)+i \phi \cdot J},
\end{align*}
$$

where

$$
\begin{equation*}
K(x, y)=-\delta^{(4)}(x-y)\left(\partial_{x}^{2}+m_{0}^{2}-i \epsilon\right) . \tag{7.82}
\end{equation*}
$$

We will need the inverse $K^{-1}$ of the linear operator $K$ defined via

$$
\begin{equation*}
\int \mathrm{d}^{4} y K(x, y) K^{-1}(y, z)=\delta^{(4)}(x-z) \tag{7.83}
\end{equation*}
$$

Inserting $K(x, y)$ yields

$$
\begin{equation*}
\left(\partial_{x}^{2}+m_{0}^{2}-i \epsilon\right) K^{-1}(x, z)=-\delta^{(4)}(x-z) . \tag{7.84}
\end{equation*}
$$

Since the Feynman propagator is a Green's function for the Klein-Gordon equation, we can just quote our result from QFT I,

$$
\begin{equation*}
K^{-1}(x, z)=-i D_{F}(x-z) \text { with } D_{F}(x)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i}{p^{2}-m_{0}^{2}+i \epsilon} e^{-i p \cdot x} \tag{7.85}
\end{equation*}
$$

Note that in this approach the $i \epsilon$-term arises automatically as a consequence of (7.80). Thus,

$$
\begin{equation*}
Z_{0}[J]=\int \mathcal{D} \phi e^{-\frac{1}{2} \phi \cdot D_{F}^{-1} \phi+i \phi \cdot J} \tag{7.86}
\end{equation*}
$$

We can manipulate this further by completing the square as

$$
\begin{equation*}
\frac{i}{2} \phi \cdot K \phi+i \phi \cdot J=\frac{i}{2}\left(\phi+K^{-1} J\right) \cdot K\left(\phi+K^{-1} J\right)-\frac{i}{2} J \cdot K^{-1} J \tag{7.87}
\end{equation*}
$$

and shift the measure $\mathcal{D} \phi=\mathcal{D} \phi^{\prime}$ with $\phi^{\prime}=\phi+K^{-1} \cdot J$ to find

$$
\begin{align*}
Z_{0}[J] & =\int \mathcal{D} \phi^{\prime} e^{\frac{i}{2} \phi^{\prime} \cdot K \phi^{\prime}} e^{-\frac{1}{2} J \cdot D_{F} J} \\
& =Z_{0}[0] e^{\frac{1}{2}(i J) \cdot D_{F}(i J)}  \tag{7.88}\\
& \equiv Z_{0}[0] e^{\frac{1}{2} \int \mathrm{~d}^{4} x \mathrm{~d}^{4} y i J(x) D_{F}(x-y) i J(y)} .
\end{align*}
$$

We notice that

$$
\begin{equation*}
Z_{0}[0]=\int \mathcal{D} \phi^{\prime} e^{-\frac{1}{2} \phi^{\prime}(-K) \phi^{\prime}} \tag{7.89}
\end{equation*}
$$

is a Gaussian integral. To compute it formally, we recall the finite dimensional formula ( $v \equiv v_{i}, i=$ $1, \ldots, n$ )

$$
\begin{equation*}
\int \mathrm{d}^{n} v e^{-\frac{1}{2} v^{T} M v}=\sqrt{\frac{(2 \pi)^{n}}{\operatorname{det} M}} \tag{7.90}
\end{equation*}
$$

for $M$ symmetric and positive definite (or more generally for $\operatorname{Re}\left(\mu_{i}\right)>0$, where $\mu_{i}$ are the eigenvalues of $M$ ). Formally this generalizes to

$$
\begin{equation*}
Z_{0}[0]=\frac{\text { const. }}{\sqrt{\operatorname{det}(-i K)}} \tag{7.91}
\end{equation*}
$$

where $\operatorname{det}(-i K)=\operatorname{det}\left[\delta^{(4)}(x-y)\left(i \partial_{x}^{2}+i m^{2}+\epsilon\right)\right]$ is called a functional determinant. Note that $Z_{0}[0]$ is a divergent quantity, but as will be discussed in the tutorials it can be defined rigorously by a regularisation procedure, e.g. on a lattice with IR and UV cutoff. For finite cutoffs, $Z_{0}[0]$ is finite. All physical quantities are computed in the regularized theory and the expression $Z_{0}[0]$ cancels in all such expressions because the Green's functions derive from $Z[J] / Z[0]$. In the end one takes the continuum limit, and the formal divergence of $\operatorname{det}(-i K)$ poses no problems.

Our result (7.88) allows us to compute the $n$-point functions

$$
\begin{equation*}
G_{0}\left(x_{1}, \ldots, x_{n}\right)=\left.\frac{\delta}{i \delta J\left(x_{1}\right)} \ldots \frac{\delta}{i \delta J\left(x_{n}\right)} \underbrace{e^{\frac{1}{2} i J \cdot D_{F} i J}}_{\equiv \frac{z_{0}[\mid]}{z_{0}[0]}}\right|_{J=0} \tag{7.92}
\end{equation*}
$$

of the free theory by differentiating the exponential with respect to $J$ and setting $J=0$. To compute a $2 n$-point function one expands the exponential precisely to $n$-th order, while the $2 n+1$-functions vanish. In this spirit the 2-point function is computed as

$$
\begin{equation*}
G_{0}\left(x_{1}, x_{2}\right)=\frac{\delta}{i \delta J\left(x_{1}\right)} \frac{\delta}{i \delta J\left(x_{2}\right)} \frac{1}{2} \int \mathrm{~d}^{4} x \mathrm{~d}^{4} y i J(x) D_{F}(x-y) i J(y)=D_{F}\left(x_{1}-x_{2}\right), \tag{7.93}
\end{equation*}
$$

confirming that the Green's function of the D'Alembert operator gives the 2-point correlator in the free theory. Feynman diagrams are a graphical way to organise the combinatorics of expanding the exponential and taking the functional derivatives. For example the result for

$$
\begin{align*}
G_{0}\left(x_{1}, x_{2}, x_{3}, x_{4}\right)= & \frac{\delta}{i \delta J\left(x_{1}\right)} \frac{\delta}{i \delta J\left(x_{2}\right)} \frac{\delta}{i \delta J\left(x_{3}\right)} \frac{\delta}{i \delta J\left(x_{4}\right)} \\
& \times \frac{1}{2!} \frac{1}{2} \int \mathrm{~d}^{4} x \mathrm{~d}^{4} y i J(x) D_{F}(x-y) i J(y) \frac{1}{2} \int \mathrm{~d}^{4} u \mathrm{~d}^{4} v i J(u) D_{F}(u-v) i J(v) \tag{7.94}
\end{align*}
$$

can be expressed graphically as

where a line between $x_{i}$ and $x_{j}$ denotes as always a factor of $D_{F}\left(x_{i}-x_{j}\right)$. Note that interactions only occur if there is a dot on crossing lines, i.e. there is no interaction in the third diagram. As one can convince oneself, the combinatorial factors work out correctly.

### 7.5 Perturbative expansion in interacting theory

We now consider an action with extra interaction terms,

$$
\begin{equation*}
S[\phi]=S_{0}[\phi]+\int \mathrm{d}^{4} x \mathcal{L}_{\mathrm{int}}(\phi) . \tag{7.95}
\end{equation*}
$$

Then

$$
\begin{align*}
Z[J] & =\int \mathcal{D} \phi e^{i S_{0}[\phi]+i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}(\phi)+i J \cdot \phi} \\
& =\int \mathcal{D} \phi e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}(\phi)} e^{i S_{0}[\phi]+i J \cdot \phi} \\
& =\int \mathcal{D} \phi e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}\left(-i \frac{\delta}{\delta j}\right)} e^{i S_{0}[\phi]+i J \cdot \phi}  \tag{7.96}\\
& =e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}\left(-i \frac{\delta}{\delta J}\right)} \int \mathcal{D} \phi e^{i S_{0}[\phi]+i J \cdot \phi} .
\end{align*}
$$

In the third line we used (7.77), or more generally

$$
\begin{equation*}
F[\phi] e^{i \phi \cdot J}=F\left[\frac{\delta}{i \delta J}\right] e^{i \phi \cdot J} \tag{7.97}
\end{equation*}
$$

for any functional $F[\phi]$. Inserting our previous result (7.88) for $\int \mathcal{D} \phi e^{i S_{0}[\phi]+i J \cdot \phi}$ gives the important result

$$
\begin{equation*}
Z[J]=e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}\left(-i \frac{\delta}{\delta j}\right)} Z_{0}[J]=Z_{0}[0] e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}\left(-i \frac{\delta}{\delta j}\right)} e^{\frac{1}{2}(i J) \cdot D_{F}(i J)} . \tag{7.98}
\end{equation*}
$$

Equivalently one can write

$$
\begin{equation*}
Z[J]=\left.Z_{0}[0] e^{\frac{1}{2} \frac{\delta}{\delta \phi} D_{F} \frac{\delta}{\delta \phi}} e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}(\phi)+i J \cdot \phi}\right|_{\phi=0} . \tag{7.99}
\end{equation*}
$$

To prove the equivalence of (7.98) and (7.99), we consider two functionals $F$ and $G$ and note that

$$
\begin{equation*}
F\left[\frac{\delta}{\delta \phi}\right] G[\phi] e^{i J \cdot \phi}=F\left[\frac{\delta}{\delta \phi}\right] G\left[\frac{\delta}{i \delta J}\right] e^{i J \cdot \phi} . \tag{7.100}
\end{equation*}
$$

Since the derivatives commute this is

$$
\begin{equation*}
G\left[\frac{\delta}{i \delta J}\right] F\left[\frac{\delta}{\delta \phi}\right] e^{i J \cdot \phi}=G\left[\frac{\delta}{i \delta J}\right] F[i J] e^{i J \cdot \phi} . \tag{7.101}
\end{equation*}
$$

In particular this implies that

$$
\begin{equation*}
G\left[\frac{\delta}{i \delta J}\right] F[i J]=\left.G\left[\frac{\delta}{i \delta J}\right] F[i J] e^{i J \cdot \phi}\right|_{\phi=0}=\left.F\left[\frac{\delta}{\delta \phi}\right] G[\phi] e^{i J \cdot \phi}\right|_{\phi=0} . \tag{7.102}
\end{equation*}
$$

The equivalence of (7.98) and (7.99) follows from this if we identify

$$
\begin{align*}
F\left[\frac{\delta}{\delta \phi}\right] & =e^{\frac{1}{2} \frac{\delta}{\delta \phi} D_{F} \frac{\delta}{\delta \phi}},  \tag{7.103}\\
G[\phi] & =e^{i \int \mathrm{~d}^{4} x \mathcal{L i n t}^{\text {in }}(\phi)} .
\end{align*}
$$

Perturbation theory can be set up from both expressions (7.98) or (7.99). Starting from (7.98) is close in spirit to our discussion of the free theory as will be demonstrated in the tutorials. In the sequel we take the approach via (7.99) and make the following definition:

$$
\begin{equation*}
\langle F[\phi]\rangle_{0}:=\left.e^{\frac{1}{2} \frac{\delta}{\delta \phi} D_{F} \frac{\delta}{\delta \phi}} F[\phi]\right|_{\phi=0} . \tag{7.104}
\end{equation*}
$$

With this notation (7.99) becomes

$$
\begin{equation*}
\frac{Z[J]}{Z[0]}=\frac{\left\langle e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}+i J \cdot \phi}\right\rangle_{0}}{\left\langle e^{\left.i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}\right\rangle_{0}}\right.} \tag{7.105}
\end{equation*}
$$

$G\left(x_{1}, \ldots, x_{n}\right)$ follows from (7.105) either by expanding $e^{i J \cdot \phi}$ and using (7.63), or directly as

$$
\begin{align*}
G\left(x_{1}, \ldots, x_{n}\right) & =\left.\frac{\delta}{i \delta J\left(x_{1}\right)} \ldots \frac{\delta}{i \delta J\left(x_{n}\right)} \frac{Z[J]}{Z[0]}\right|_{J=0} \\
& =\left.\frac{\left\langle\frac{\delta}{i \delta J\left(x_{1}\right)} \cdots \frac{\delta}{i \delta J\left(x_{n}\right)} e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}+i J \cdot \phi}\right\rangle_{0}}{\left\langle e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}}\right\rangle_{0}}\right|_{J=0} \tag{7.106}
\end{align*}
$$

We can perform the functional derivative in the nominator and arrive at the final expression

$$
\begin{equation*}
G\left(x_{1}, \ldots, x_{n}\right)=\frac{\left\langle\phi\left(x_{1}\right) \ldots \phi\left(x_{n}\right) e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}}\right\rangle_{0}}{\left\langle e^{i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}}\right\rangle_{0}} \tag{7.107}
\end{equation*}
$$

To understand the meaning of this expression we must familiarize ourselves more with the operation (7.104). First, let us compute

$$
\begin{equation*}
\left\langle\phi\left(x_{1}\right) \phi\left(x_{2}\right)\right\rangle_{0}:=\left.e^{\frac{1}{2} \frac{\delta}{2 \phi} D_{F} \frac{\delta}{\delta \phi}} \phi\left(x_{1}\right) \phi\left(x_{2}\right)\right|_{\phi=0} . \tag{7.108}
\end{equation*}
$$

Expanding the exponential and keeping precisely terms of quadratic order in $\frac{\delta}{\delta \phi}$ yields

$$
\begin{equation*}
\frac { 1 } { 2 } \int \mathrm { d } ^ { 4 } x \mathrm { d } ^ { 4 } y \frac { \delta } { \delta \phi ( x ) } D _ { F } ( x - y ) \frac { \delta } { \delta \phi ( y ) } \phi ( x _ { 1 } ) \phi ( x _ { 2 } ) | _ { \phi = 0 } = D _ { F } ( x _ { 1 } - x _ { 2 } ) \equiv \longdiv { \phi ( x _ { 1 } ) \phi ( x _ { 2 } ) } . \tag{7.109}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\left\langle\phi\left(x_{1}\right) \phi\left(x_{2}\right)\right\rangle_{0}=\langle 0| T \hat{\phi}\left(x_{1}\right) \hat{\phi}\left(x_{2}\right)|0\rangle, \tag{7.110}
\end{equation*}
$$

where $|0\rangle$ is the free vacuum. Generalising this reasoning gives a simple and very direct proof of Wick's theorem, which will be presented in the exercises. We will end up with the result known from QFT I:


- $\left\langle\phi_{1} \ldots \phi_{2 n+1}\right\rangle_{0}=0$.

With Wick's theorem at hand we go back to $G\left(x_{1}, \ldots x_{n}\right)$ and $\operatorname{expand} \exp \left[i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}\right]$ in (7.107) to recover exactly the same Feynman rules as in operator language from the Gell-Mann-Low formula in the interaction picture. We refer to QFT I for a precise statement of these rules in position and in momentum space. The special case $\exp \left[i \int \mathrm{~d}^{4} x \mathcal{L}_{\text {int }}\right]=1$ gives the free theory.
Note that $\phi\left(x_{i}\right)$ in (7.107) or $\frac{\delta}{i \delta J\left(x_{i}\right)}$ in (7.106) corresponds to an external point $x_{i}$ in the Feynman diagram. Thus $Z[0]$ contains no external points and therefore represents the partition function since

$$
\begin{equation*}
Z[0]=e^{\Sigma_{i} V_{i}} . \tag{7.111}
\end{equation*}
$$

Here $V_{i}$ is the value of a vacuum bubble labeled by $i$. Consequently $\frac{Z[J]}{Z[0]}$, i.e. the generating functional of the Green's functions, contains no vacuum bubbles.

## On the counting of loops and factors of $\hbar$

Consider a fully connected Feynman diagram in momentum space with

- $E=$ number of external lines,
- $I=$ number of internal lines,
- $V=$ number of vertices and
- $L=$ number of loops.
$L$ is the number of unfixed momentum integrals. From the momentum space Feynman rules we recall that
- each internal propagator gives one momentum integral,
- each vertex gives one $\delta^{(4)}\left(\Sigma_{i} p_{i}\right)$, but
- one linear combination of these delta-functions merely corresponds to overall momentum conservation.

This proves Euler's formula

$$
\begin{equation*}
L=I-V+1 \text {. } \tag{7.112}
\end{equation*}
$$

The perturbative expansion in loops corresponds to an expansion in $\hbar$. To see this, we note that the exponent in the path integral is - reinstating $\hbar$ -

$$
\begin{equation*}
\frac{i}{\hbar} S=\frac{i}{\hbar} \int \mathrm{~d}^{4} x \phi\left(-\frac{1}{2}\left(\partial^{2}+m^{2}\right) \phi\right)+\frac{i}{\hbar} \int \mathrm{~d}^{4} x \mathcal{L}_{\mathrm{int}} . \tag{7.113}
\end{equation*}
$$

Each internal line carries a propagator, i.e. a factor of

$$
\begin{equation*}
\left(\frac{\partial^{2}+m^{2}}{\hbar}\right)^{-1} \sim \hbar \tag{7.114}
\end{equation*}
$$

and each vertex carries a factor of $\hbar^{-1}$. For a given scattering process with fixed number $E$ of external lines, expansion of a scattering amplitude in $L$ therefore gives $I-V=L-1$ extra powers of $\hbar$ (which come in addition to factors of $\hbar$ from the external particles, whose power is fixed for given $E$ ). Taking into account these extra factors, in most field theories the classical scattering amplitude, i.e. $O\left(\hbar^{0}\right)$, corresponds to the tree-level diagram, but exceptions to this rule of thumb exist. ${ }^{5}$

Note furthermore that for $\mathcal{L}_{\text {int }}=\frac{\lambda}{n!} \phi^{n}$, a loop expansion with fixed $E$ is an expansion in $\lambda$, because

$$
\begin{equation*}
E+2 I=n V \tag{7.115}
\end{equation*}
$$

since both ends of each internal line and one end of each external line end in an n-point vertex. With $L=I-V+1$ this gives

$$
\begin{equation*}
(n-2) V=2 L+(E-2) . \tag{7.116}
\end{equation*}
$$

So for $E$ fixed, an expansion in $L$ corresponds to an expansion in $V$.

### 7.6 The Schwinger-Dyson equation

One advantage of the path-integral quantisation method over the canonical formalism is that considerations of symmetry are more transparent. In particular it becomes very clear that classical symmetries carry over to the quantum theory - but only provided the path integral measure is invariant.

To see this we consider a general theory with arbitrary fields collectively denoted by $\phi(x)$ defined by

$$
\begin{equation*}
Z[J]=\int \mathcal{D} \phi e^{i(S[\phi]+\phi \cdot J)} \tag{7.117}
\end{equation*}
$$

Assume that a transformation

$$
\begin{equation*}
\phi(x) \rightarrow \phi^{\prime}(x)=\phi(x)+\epsilon \Delta \phi(x)+O\left(\epsilon^{2}\right) \tag{7.118}
\end{equation*}
$$

leaves the measure $\mathcal{D} \phi$ invariant, i.e.

$$
\begin{equation*}
\mathcal{D} \phi=\mathcal{D} \phi^{\prime} \tag{7.119}
\end{equation*}
$$

Note that this is really an assumption which may fail in general. Then

$$
\begin{align*}
Z[J] & =\int \mathcal{D} \phi e^{i(S[\phi]+\phi \cdot J)} \\
& =\int \mathcal{D} \phi^{\prime} e^{i\left(S\left[\phi^{\prime}\right]+\phi^{\prime} \cdot J\right)} \\
& =\int \mathcal{D} \phi e^{i\left(S\left[\phi^{\prime}\right]+\phi^{\prime} \cdot J\right)}  \tag{7.120}\\
& =\int \mathcal{D} \phi e^{i(S[\phi]+\phi \cdot J)}\left[1+i\left(\frac{\delta S}{\delta \phi}+J\right) \epsilon \Delta \phi+O\left(\epsilon^{2}\right)\right]
\end{align*}
$$

[^35]In the second line we simply relabeled the integration variable and in the third we used invariance of the measure. Thus

$$
\begin{align*}
0 & =\int \mathcal{D} \phi e^{i(S[\phi]+\phi \cdot J)}\left[\frac{\delta S}{\delta \phi} \Delta \phi+J \Delta \phi\right] \\
& \equiv \int \mathcal{D} \phi e^{i(S[\phi]+\phi \cdot J)}\left[\int \mathrm{d}^{4} y\left(\frac{\delta S}{\delta \phi(y)}+J(y)\right) \Delta \phi(y)\right] . \tag{7.121}
\end{align*}
$$

For the special transformation $\Delta \phi(y)=\delta^{(4)}(y-x)$, our assumption (7.119) does indeed hold, and therefore

$$
\begin{equation*}
\int \mathcal{D} \phi\left(\frac{\delta S}{\delta \phi(x)}+J(x)\right) e^{i(S[\phi]+J \phi)}=0 \tag{7.122}
\end{equation*}
$$

Equ. (7.122) is the celebrated Schwinger-Dyson equation. It states that the classical equation of motion (here in presence of a source $J$ ),

$$
\begin{equation*}
0=\frac{\delta S}{\delta \phi}+J \tag{7.123}
\end{equation*}
$$

holds as an operator equation in the quantum theory, i.e. it holds inside the path integral (with no further field insertion at the same spacetime point).
If we act on (7.122) with $\prod_{i=1}^{n} \frac{\delta}{i \delta J\left(x_{i}\right)}$ we arrive at

$$
\begin{align*}
\int \mathcal{D} \phi \quad & {\left[\left(\frac{\delta S}{\delta \phi(x)}+J(x)\right) \phi\left(x_{1}\right) \ldots \phi\left(x_{n}\right)+\sum_{i=1}^{n} \phi\left(x_{1}\right) \ldots \phi\left(x_{i-1}\right)\left(-i \delta\left(x-x_{i}\right)\right) \phi\left(x_{i+1}\right) \ldots \phi\left(x_{n}\right)\right] } \\
& \times e^{i(S[\phi]+J \phi)}=0 . \tag{7.124}
\end{align*}
$$

Taking $J=0$ gives

$$
\begin{equation*}
\langle\Omega| T \frac{\delta S}{\delta \phi(x)} \prod_{i=1}^{n} \phi\left(x_{i}\right)|\Omega\rangle=i \sum_{i=1}^{n}\langle\Omega| T \phi\left(x_{1}\right) \phi\left(x_{i-1}\right) \delta\left(x-x_{i}\right) \phi\left(x_{i+1}\right) \ldots \phi\left(x_{n}\right)|\Omega\rangle . \tag{7.125}
\end{equation*}
$$

In this form the Schwinger-Dyson equation shows that the equations of motion hold inside a correlation function up to contact terms on the right-hand side.

As a simple corollary we derive the Ward-Takahashi identities: To this end start from (7.121) again and consider a transformation

$$
\begin{equation*}
\phi \rightarrow \phi+\epsilon \Delta \phi \tag{7.126}
\end{equation*}
$$

with

$$
\begin{equation*}
\Delta \phi(y)=\delta^{(4)}(y-x) \delta \phi(x) \tag{7.127}
\end{equation*}
$$

for a continuous global classical symmetry $\delta \phi$ such that

$$
\begin{equation*}
\frac{\delta S}{\delta \phi(x)} \delta \phi(x)=-\partial_{\mu} j^{\mu}(x) \tag{7.128}
\end{equation*}
$$

(with no integration over the $x$-variable on the left) holds off-shell. Here $j^{\mu}(x)$ is the Noether current associated with the global continuous symmetry. If this transformation leaves the measure invariant, then (7.121) yields

$$
\begin{equation*}
0=\int \mathcal{D} e^{i(S[\phi]+\phi J)}\left(-\partial_{\mu} j^{\mu}(x)+J(x) \delta \phi(x)\right) . \tag{7.129}
\end{equation*}
$$

Taking $\prod_{i=1}^{n} \frac{\delta}{i \delta J\left(x_{i}\right)}$ at $J=0$ gives the Ward-Takahashi identity:

$$
\begin{equation*}
\partial_{\mu}\langle\Omega| T j^{\mu} \prod_{i=1}^{n} \phi\left(x_{i}\right)|\Omega\rangle=-i \sum_{i=1}^{n}\langle\Omega| T \phi\left(x_{1}\right) \ldots \phi\left(x_{i-1}\right) \delta \phi(x) \delta^{(4)}\left(x-x_{i}\right) \phi\left(x_{i+1}\right) \ldots \phi\left(x_{n}\right)|\Omega\rangle . \tag{7.130}
\end{equation*}
$$

While we have seen this equation already in QFT1, the above path integral derivation indeed clarifies, as promised, that the Ward identities only hold for symmetries that leave the measure invariant. If this is not the case, the symmetry is said to be anomalous and current conservation (up to contact terms) does not hold at the quantum level. We will learn more about anomalies later in this course.

To conclude this chapter we go back to the Schwinger-Dyson equation in its form (7.122) and use the identity

$$
\begin{equation*}
F[\phi] e^{i \phi \cdot J}=F\left[\frac{\delta}{i \delta J}\right] e^{i \phi \cdot J} \tag{7.131}
\end{equation*}
$$

for $F[\phi]=\frac{\delta S}{\delta \phi}$ to rewrite (7.122) as

$$
\begin{equation*}
\left(\left.\frac{\delta S}{\delta \phi}\right|_{\phi=\frac{\delta}{i \delta J}}+J\right) Z[J]=0 \tag{7.132}
\end{equation*}
$$

The Schwinger-Dyson equation in this form is a fundamental identity for $Z[J]$ that beautifully encodes the structure of the correlation functions. To see this we introduce the abbreviations

$$
\begin{equation*}
\phi\left(x_{i}\right) \equiv \phi_{i}, \quad J\left(x_{i}\right) \equiv J_{i} \tag{7.133}
\end{equation*}
$$

and

$$
\begin{equation*}
S[\phi]=-\frac{1}{2} \phi_{i} D_{F}^{-1}{ }_{i j} \phi_{j}+\sum_{m \geq 3} \frac{1}{m!} Y_{i_{n} \ldots i_{m}}^{(m)} \phi_{i_{1}} \ldots \phi_{i_{m}} \tag{7.134}
\end{equation*}
$$

as shorthand for

$$
\begin{equation*}
-\frac{1}{2} \int \mathrm{~d}^{4} x_{i} \mathrm{~d}^{4} x_{j} \phi\left(x_{i}\right) \delta^{(4)}\left(x_{i}-x_{j}\right)\left(\partial_{i}^{2}+m_{0}^{2}\right) \phi\left(x_{j}\right) \tag{7.135}
\end{equation*}
$$

and similarly for the interaction terms. To evaluate (7.132), i.e.

$$
\begin{equation*}
\left.\frac{\delta S}{\delta \phi_{i}}\right|_{\phi_{i}=\frac{\delta}{\delta \delta_{i}}} Z[J]=-J_{i} Z[J], \tag{7.136}
\end{equation*}
$$

we note that

$$
\begin{equation*}
\left.\frac{\delta S}{\delta \phi_{i}}\right|_{\phi_{i}=\frac{\delta}{i \delta J_{i}}}=-D_{F}^{-1}{ }_{i j} \frac{\delta}{i \delta J_{j}}+\sum_{m} \frac{1}{(m-1)!} Y_{i i_{1} \ldots i_{m-1}}^{m} \frac{\delta}{i \delta J_{i_{1}}} \ldots \frac{\delta}{i \delta J_{i_{m-1}}} . \tag{7.137}
\end{equation*}
$$

Now multiply both sides by $\left(D_{F}\right)_{k i}$ using $\left(D_{F}\right)_{k i}\left(D_{F}^{-1}\right)_{i j}=\delta_{k j}$ to find

$$
\begin{equation*}
-\frac{\delta}{i \delta J_{k}} Z[J]+\left(D_{F}\right)_{k i} \sum_{m \geq 3} \frac{1}{(m-1)!} Y_{i_{1} \ldots i_{m-1}}^{(m)} \frac{\delta}{i \delta J_{i_{1}}} \ldots \frac{\delta}{i \delta J_{i_{m-1}}} Z[J]=-\left(D_{F}\right)_{k i} J_{i} Z[J] . \tag{7.138}
\end{equation*}
$$

Thus the Schwinger-Dyson equation can be written as

$$
\begin{equation*}
\frac{\delta}{i \delta J_{k}} Z[J]=\left(D_{F}\right)_{k i}\left(J_{i}+\sum_{m} \frac{1}{(m-1)!} Y_{i i_{1} \ldots i_{m-1}}^{m} \frac{\delta}{i \delta J_{i_{1}}} \ldots \frac{\delta}{i \delta J_{i_{m-1}}}\right) Z[J] . \tag{7.139}
\end{equation*}
$$

Let us interpret this graphically by making the replacements

$$
\begin{array}{r}
Z[J] \leftrightarrow(Z) \quad \prod_{i=1}^{n} \frac{\delta}{i \delta J_{i}} Z[J]=\left.Z\right|_{n} \quad(n \text { lines attached }), \\
J_{i} \leftrightarrow x_{i}, \quad\left(D_{F}\right)_{i j}=i-j, \quad Y_{i_{n} \ldots i_{m}}=i_{m}, \ldots
\end{array}
$$

Thus the Schwinger-Dyson equation has the graphic representation


But this just states how the full 1-point function can be understood in terms of full higher $n$-point functions and bare interactions as encoded in $S[\phi]$ : The lefthand side gives the amplitude for a particle starting at $x_{k}$ to couple to the vacuum via a full quantum process, represented by (Z. According to the righthand-side, it can either be absorbed by a classical source $J$ at $x_{i}$, where we must sum over all possible insertion locations $x_{i}$, or it can decay first into two or more particles by a tree-level process before these in turn couple to the full quantum process. The numerical factors are the appropriate symmetry factors.

In fact, this interpretation of the Schwinger-Dyson equation suggests a very elegant derivation of the key concepts of Quantum Field Theory in a manner completely reverse to the approach we have taken in this course: Start with the concept of a free particle as an object that propagates freely in spacetime, and define the propagator $D_{F i j}$ as the quantum probability amplitude for a free particle to travel from $i$ to $j$. Next introduce the action $S[\phi]$ as in equ. (7.134) as nothing but a symbolic means to encode all perturbative bare interactions $Y_{i_{1} \ldots i_{m}}$, i.e. all the ways how an interacting particle can split up into other particles along its way. Then the Schwinger-Dyson equation in its form (7.140) is an obvious property of quantum probability amplitudes in the sense described in the previous paragraph. Taking (7.140) as a starting point and going backwards leads to (7.132) as an equation for $Z[J]$, which can be solved by making the Fourier ansatz $Z[J]=\int \mathcal{D} \phi \tilde{Z}[\phi] e^{i \phi \cdot J}$ and solving the corresponding equation for the Fourier transform $\tilde{Z}[\phi]$ with the result

$$
\begin{equation*}
Z[J] \sim \int \mathcal{D} \phi e^{i(S[\phi]+\phi \cdot J)} . \tag{7.141}
\end{equation*}
$$

More details on this approach can be found in the QFT notes by Cvitanovic and the book by Banks.

### 7.7 Connected diagrams

As a consequence of the denominator in (7.107), a general Green's function $G\left(x_{1}, \ldots, x_{n}\right)$ is computed only from those Feynman diagrams which contain no vacuum bubbles, i.e. no diagrams not connected to any of the external points. Nonetheless, $G\left(x_{1}, \ldots, x_{n}\right)$ does receive contributions from partially connected Feynman diagrams. These are diagrams that factor into subdiagrams each of which is connected only to some of the $n$ external points $x_{1}, \ldots, x_{n}$. As established by the LSZ formalism, what enters the computation of scattering amplitudes are only the fully connected Green's functions $G^{(c)}\left(x_{1}, \ldots, x_{n}\right)$ corresponding to fully connected Feynman graphs, i.e. to those Feynman diagrams which do not factor into subdiagrams.
Let us denote by the effective action $i W[J]$ the generating functional for all fully connected Green's functions $G^{(c)}\left(x_{1},, x_{n}\right),{ }^{6}$

$$
\begin{equation*}
i W[J]:=\sum_{n=1}^{\infty} \frac{1}{n!} \int \mathrm{d}^{4} x_{1} \ldots \int \mathrm{~d}^{4} x_{n}\left(i J\left(x_{1}\right)\right) \ldots\left(i J\left(x_{n}\right)\right) G^{(c)}\left(x_{1}, \ldots, x_{n}\right), \tag{7.142}
\end{equation*}
$$

to be contrasted with the generating functional for all Green's functions $G\left(x_{1},, x_{n}\right)$,

$$
\begin{equation*}
\frac{Z[J]}{Z[0]}=\sum_{n=0}^{\infty} \frac{1}{n!} \int \mathrm{d}^{4} x_{1} \ldots \int \mathrm{~d}^{4} x_{n}\left(i J\left(x_{1}\right)\right)\left(i J\left(x_{n}\right)\right) G\left(x_{1}, \ldots, x_{n}\right) . \tag{7.143}
\end{equation*}
$$

One can convince oneself that

$$
\begin{equation*}
\frac{Z[J]}{Z[0]}=1+i W[J]+\frac{1}{2}(i W[J])^{2}+\frac{1}{3!}(i W[J])^{3}+\ldots \tag{7.144}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\frac{Z[J]}{Z[0]}=e^{i W[J]} \text {, i.e. } i W[J]=\ln \frac{Z[J]}{Z[0]} . \tag{7.145}
\end{equation*}
$$

To see (7.144) one observes that forming products of all fully connected Green's functions indeed gives all (potentially partially connected) ones. The factor $1 / n!$ prevents overcounting and is necessary to match the numerical factors in (7.143). For later purposes we define

$$
\begin{equation*}
\tau\left(x_{1}, \ldots x_{n}\right)=\frac{\delta}{i \delta J\left(x_{1}\right)} \ldots \frac{\delta}{i \delta J\left(x_{n}\right)} i W[J] \tag{7.146}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\left.\tau\left(x_{1}, \ldots x_{n}\right)\right|_{J=0}=G^{(c)}\left(x_{1}, \ldots, x_{n}\right) . \tag{7.147}
\end{equation*}
$$

[^36]
### 7.8 The 1PI effective action

An important subclass of fully connected Feynman diagrams are the 1-particle-irreducible (1PI) diagrams, which cannot be cut into two non-trivial diagrams by cutting a single internal line. We would like to establish a generating functional for the 1PI connected diagrams, called 1PI effective action $\Gamma[\varphi]$. Let us first define $\Gamma[\varphi]$ and then discuss its relation to the 1PI diagrams.

- If $\phi(x)$ denotes the field in $\mathcal{L}(\phi)$, define

$$
\begin{equation*}
\varphi_{J}(x):=\frac{\delta W[J]}{\delta J(x)} . \tag{7.148}
\end{equation*}
$$

By definition this is the 1-point function of $\phi(x)$ in the presence of the source $J$,

$$
\begin{equation*}
\varphi_{J}(x)=\langle\Omega| \hat{\phi}(x)|\Omega\rangle_{J}, \tag{7.149}
\end{equation*}
$$

as will be shown in the tutorials. The subscript in $\varphi_{J}(x)$ stresses the dependence on $J$. Then

$$
\begin{equation*}
\left.\varphi_{J}(x)\right|_{J=0}=\langle\Omega| \hat{\phi}(x)|\Omega\rangle . \tag{7.150}
\end{equation*}
$$

Without loss of generality we assume that $\left.\varphi_{J}(x)\right|_{J=0}=0$ : By Lorentz invariance $\langle\Omega| \hat{\phi}(x)|\Omega\rangle=$ $\langle\Omega| \hat{\phi}(0)|\Omega\rangle=c \forall x$, and if $c \neq 0$ we redefine $\phi(x) \rightarrow \phi(x)-c$ to ensure $\left.\varphi_{J}(x)\right|_{J=0}=0$.

- Assume that the relation between $\varphi_{J}$ and $J$ is invertible such that for a given function $\varphi(x)$ there exists a unique $J_{\varphi}(x)$ with the property

$$
\begin{equation*}
\langle\Omega| \hat{\phi}(x)|\Omega\rangle_{J}=\varphi(x) . \tag{7.151}
\end{equation*}
$$

Then we define the 1PI effective action $\Gamma[\varphi]$ to be the Legendre transform of $W[J]$, i.e.

$$
\begin{equation*}
\Gamma[\varphi]:=W\left[J_{\varphi}\right]-\varphi \cdot J_{\varphi} \equiv W\left[J_{\varphi}\right]-\int \mathrm{d}^{4} x^{\prime} \varphi\left(x^{\prime}\right) J_{\varphi}\left(x^{\prime}\right) . \tag{7.152}
\end{equation*}
$$

$\Gamma[\varphi]$ is a functional of $\varphi$ only, i.e.

$$
\begin{equation*}
\delta \Gamma[\varphi]=\frac{\delta \Gamma}{\delta \varphi} \cdot \delta \varphi \tag{7.153}
\end{equation*}
$$

with

$$
\begin{align*}
\frac{\delta \Gamma}{\delta \varphi(x)} & =\frac{\delta W}{\delta \varphi(x)}-\frac{\delta \varphi}{\delta \varphi(x)} \cdot J_{\varphi}-\varphi \cdot \frac{\delta J_{\varphi}}{\delta \varphi(x)} \\
& =\int \mathrm{d}^{4} y \underbrace{\frac{\delta W}{\delta J(y)}}_{=\varphi} \frac{\delta J(y)}{\delta \varphi(x)}-J_{\varphi}(x)-\varphi \cdot \frac{\delta J}{\delta \varphi(x)} \tag{7.154}
\end{align*}
$$

and thus

$$
\begin{equation*}
\frac{\delta \Gamma}{\delta \varphi(x)}=-J_{\varphi}(x) \tag{7.155}
\end{equation*}
$$

## A note on the assumption of invertibility

In the Euclidean theory (after Wick rotation) one can show that $W_{E}[J]$ is convex, which implies that $J_{\varphi}$ is well-defined as a function of $\varphi^{7}$ Thus $\Gamma_{E}[\varphi]$ is well-defined as the Legendre trans-

[^37]formation of $W_{E}[J]$ in the Euclidean theory. All computations can be performed there and the final results can be Wick rotated back in the end.

Our claim is now that $\Gamma[\varphi]$ is the generating functional for the 1PI connected amputated Green's functions $\Gamma_{n}\left(x_{1}, \ldots, x_{n}\right)$,

$$
\begin{equation*}
i \Gamma[\varphi]=\sum_{n=0}^{\infty} \frac{1}{n!} \int \mathrm{d}^{4} x_{1} \ldots \mathrm{~d}^{4} x_{n} \varphi\left(x_{1}\right) \varphi\left(x_{n}\right) \Gamma_{n}\left(x_{1}, x_{n}\right) \tag{7.156}
\end{equation*}
$$

Put differently, if we define

$$
\begin{equation*}
\tilde{\Gamma}_{n}\left(x_{1}, \ldots, x_{n}\right):=\frac{\delta}{\delta \varphi\left(x_{1}\right)} \ldots \frac{\delta}{\delta \varphi\left(x_{n}\right)} i \Gamma[\varphi], \tag{7.157}
\end{equation*}
$$

then the 1PI connected amputated Green's functions $\Gamma_{n}\left(x_{1}, x_{n}\right)$ are

$$
\begin{equation*}
\Gamma_{n}\left(x_{1}, \ldots, x_{n}\right)=\left.\tilde{\Gamma}_{n}\left(x_{1}, \ldots, x_{n}\right)\right|_{\varphi=0} \tag{7.158}
\end{equation*}
$$

To prove this, we establish a relation between

$$
\begin{equation*}
\tau_{n}\left(x_{1}, \ldots, x_{n}\right)=\frac{\delta}{i \delta J\left(x_{1}\right)} \ldots \frac{\delta}{i \delta J\left(x_{n}\right)} i W[J] \tag{7.159}
\end{equation*}
$$

and $\tilde{\Gamma}_{n}\left(x_{1}, \ldots, x_{n}\right)$ that justifies this claim.

1) $n=2$ :

We have

$$
\begin{align*}
\tilde{\Gamma}_{2}\left(x_{1}, x_{2}\right) & =\frac{\delta}{\delta \varphi\left(x_{1}\right)} \frac{\delta}{\delta \varphi\left(x_{2}\right)} i \Gamma[\varphi]=\frac{\delta}{\delta \varphi\left(x_{1}\right)}\left(-i J\left(x_{2}\right)\right), \\
\tau_{2}\left(x_{1}, x_{2}\right) & =\frac{\delta}{i \delta J\left(x_{1}\right)} \frac{\delta}{i \delta J\left(x_{2}\right)} i W[J]=\frac{\delta}{i \delta J\left(x_{1}\right)} \varphi\left(x_{2}\right) \tag{7.160}
\end{align*}
$$

and thus using the chain rule

$$
\begin{align*}
\int \mathrm{d} x_{2} \tau_{2}\left(x_{1}, x_{2}\right) \tilde{\Gamma}_{2}\left(x_{2}, x_{3}\right) & =\int \mathrm{d} x_{2} \frac{\delta \varphi\left(x_{2}\right)}{i \delta J\left(x_{1}\right)} \frac{\delta}{\delta \varphi\left(x_{2}\right)}\left(-i J\left(x_{3}\right)\right) \\
& =-\frac{\delta J\left(x_{3}\right)}{i \delta J\left(x_{1}\right)}=-\delta\left(x_{1}-x_{3}\right) . \tag{7.161}
\end{align*}
$$

Therefore

$$
\begin{equation*}
\tau_{2}\left(x_{1}, x_{2}\right)=-\tilde{\Gamma}_{2}^{-1}\left(x_{1}, x_{2}\right) \tag{7.162}
\end{equation*}
$$

Evaluated at $J=0$ and thus also $\varphi=0\left(\right.$ since $\left.\left.\varphi\right|_{J=0}=0\right)$ this gives

$$
\begin{equation*}
G_{2}^{(c)}\left(x_{1}, x_{2}\right)=-\Gamma_{2}^{-1}\left(x_{1}, x_{2}\right) \tag{7.163}
\end{equation*}
$$

We will come back to the interpretation of this relation after the treating the case $n=3$.

## 2) $n=3$

To proceed further it is useful to systematise our computations as follows: First by means of the chain rule,

$$
\begin{align*}
\frac{\delta}{i \delta J(x)} & =\int \mathrm{d}^{4} y \frac{\delta \varphi(y)}{i \delta J(x)} \frac{\delta}{\delta \varphi(y)}=\int \mathrm{d}^{4} y \frac{\delta^{2} i W}{i \delta J(x) i \delta J(y)} \frac{\delta}{\delta \varphi(y)}  \tag{7.164}\\
& =\int \mathrm{d}^{4} y \tau_{2}(x, y) \frac{\delta}{\delta \varphi(y)}
\end{align*}
$$

Next, we introduce the graphical representation

$$
\begin{equation*}
\tilde{\Gamma}\left(x_{1}, \ldots, x_{n}\right) \equiv\left(\left.\right|_{n},\left.\quad \tau_{n}\left(x_{1}, \ldots, x_{n}\right) \equiv(W)\right|_{n}\right. \tag{7.165}
\end{equation*}
$$

with $\mathbf{n}$ legs (denoted as $\left.\right|_{n}$ ) respectively. For example our result $\tau_{2}(x, y)=-\tilde{\Gamma}_{2}(x, y)^{-1}$ is graphically represented as

$$
\begin{equation*}
x-W-y=-(x-\Gamma-y)^{-1} \tag{7.166}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\delta}{i \delta J(x)} \equiv x-W \underbrace{-}_{\text {Integrate over } \mathrm{y}} \frac{\delta}{\delta \varphi(y)} . \tag{7.167}
\end{equation*}
$$

Finally, by definition $\frac{\delta}{i \delta J(x)}$ adds a further line to $\tau_{n}\left(x_{1}, \ldots x_{n}\right)$ and $\frac{\delta}{\delta \varphi(x)}$ adds another line to $\tilde{\Gamma}_{n}\left(x_{1}, \ldots, x_{n}\right)$. Now we compute graphically

$$
\begin{align*}
\tau_{3}(x, y, z) & \equiv \frac{\delta}{i \delta J(x)} y-W-z \\
& =x-W-\omega \frac{\delta}{\delta \varphi(\omega)} y-W-z  \tag{7.168}\\
& =x-W-\omega \frac{\delta}{\delta \varphi(\omega)}(-y-\Gamma-z)^{-1} .
\end{align*}
$$

To compute $\frac{\delta}{\delta \phi(\omega)} \tilde{\Gamma}_{2}(x, z)^{-1}$ we generalize the identity

$$
\begin{equation*}
\left(\frac{\mathrm{d}}{\mathrm{~d} \lambda} M(\lambda)^{-1}\right)_{A B}=-M_{A C}^{-1}\left(\frac{\mathrm{~d}}{\mathrm{~d} \lambda} M(\lambda)\right)_{C D} M_{D B}^{-1} \tag{7.169}
\end{equation*}
$$

for matrices $M(\lambda)$ to functional operators,

$$
\begin{align*}
\frac{\delta}{\delta \phi(\omega)} \tilde{\Gamma}_{2}^{-1}(y, z) & =-\int \mathrm{d}^{4} u \mathrm{~d}^{4} v \underbrace{\tilde{\Gamma}_{2}^{-1}(y, u)}_{=-\tau_{2}} \underbrace{\frac{\delta}{\delta \phi(\omega)} \tilde{\Gamma}_{2}(u, v)}_{=\tilde{\Gamma}_{3}(\omega, u, v)} \underbrace{\tilde{\Gamma}_{2}^{-1}(v, z)}_{=-\tau_{2}}  \tag{7.170}\\
& =-\int \mathrm{d}^{4} u \mathrm{~d}^{4} v \tau_{2}(y, u) \tilde{\Gamma}_{3}(\omega, u, v) \tau_{2}(v, z) .
\end{align*}
$$

Thus

$$
\begin{equation*}
\frac{\delta}{\delta \varphi(\omega)} y-W-z=y-W-\Gamma_{\omega}-W-z, \tag{7.171}
\end{equation*}
$$

where the subscript indicates an extra leg. Putting everything together this just means


Here full circles indicate (W) , while empty circles indicate (Г). Thus $\Gamma_{3}$ gives indeed the amputated 1PI connected 3-point function, as external legs merely carry corrections due to the fully resummed propagator

$$
\begin{equation*}
G_{2}=-W-\left.\right|_{J=0}= \tag{7.172}
\end{equation*}
$$

This can be generalised (grapically) to higher $n$ and the claim (7.156) is proven inductively for all $\Gamma_{n}\left(x_{1}, \ldots, x_{n}\right)$ with $n \geq 3$.

Let us come back to the interpretation of $\Gamma_{2}\left(x_{1}, x_{2}\right)$, which is a bit special in the following sense: Namely, the difference between $\Gamma_{2}\left(x_{1}, x_{2}\right)$ and $\Gamma_{n}\left(x_{i}\right), n \geq 3$ is that $\Gamma_{n}\left(x_{i}\right), n \geq 3$ is the sum of the tree-level vertex plus of the 1- and higher-loop 1PI amputed corrections, while $\Gamma_{2}\left(x_{1}, x_{2}\right)$ consists of the 1- and higher loop 1PI amputated diagrams minus the tree-level propagator.
This can be seen by noting that relation (7.163), i.e. $G_{2}^{(c)} \cdot \Gamma_{2}=-\mathbb{1}$, reads diagrammatically

$$
\begin{equation*}
x-W-y=-x-W-\text { - }-W-W-y \text {. } \tag{7.173}
\end{equation*}
$$

Up to an important overall minus sign this is of the same form as the identity for $n=3$ and higher $n$. Note that evaluating (7.173) at tree-level gives


We can then attach two external lines to $\Gamma_{2}\left(x_{1}, x_{2}\right)$ and consider the resulting object


To identify $\left.\Gamma_{2}\right|_{\text {loop }}$ we go to Fourier space. Let us define $\Gamma_{2}\left(p^{2}\right)$ as

$$
\begin{equation*}
\Gamma_{2}\left(x_{1}, x_{2}\right)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \Gamma_{2}\left(p^{2}\right) e^{-i p\left(x_{1}-x_{2}\right)} \tag{7.176}
\end{equation*}
$$

In momentum space, by factoring out the external lines from the diagrams, we see that

$$
\begin{align*}
& -x_{1}-x_{2} \Leftrightarrow-\frac{i}{p^{2}-m_{0}^{2}} \frac{p^{2}-m_{0}^{2}}{i} \frac{i}{p^{2}-m_{0}^{2}} \\
& x_{1}--i M^{2}-x_{2} \Leftrightarrow \frac{i}{p^{2}-m_{0}^{2}}\left(-\left.i \Gamma_{2}\left(p^{2}\right)\right|_{\mathrm{loop}}\right) \frac{i}{p^{2}-m_{0}^{2}} \tag{7.177}
\end{align*}
$$

and thus

$$
\begin{equation*}
\Gamma_{2}\left(p^{2}\right)=-\frac{p^{2}-m_{0}^{2}}{i}-\left.i \Gamma_{2}\left(p^{2}\right)\right|_{\text {loop }} \tag{7.178}
\end{equation*}
$$

On the other hand, consider the Fourier transform of the fully resummed Feynman propagator $G_{2}^{(c)}\left(x_{1}, x_{2}\right)$ via

$$
\begin{equation*}
G_{2}^{(c)}\left(x_{1}, x_{2}\right)=-\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} G_{2}^{(c)}\left(p^{2}\right) e^{-i p\left(x_{1}-x_{2}\right)}, \tag{7.179}
\end{equation*}
$$

where $G_{2}^{(c)}\left(p^{2}\right)$ is given by Dyson resummation as

$$
\begin{equation*}
G_{2}^{(c)}\left(p^{2}\right)=\frac{i}{p^{2}-m_{0}^{2}-M^{2}\left(p^{2}\right)} . \tag{7.180}
\end{equation*}
$$

Here $-i M^{2}\left(p^{2}\right)$ is the value of all amputated 1PI diagrams (i.e. without external legs) evaluated at 1-loop and higher. Relation (7.163), i.e. $G_{2}^{(c)} \cdot \Gamma_{2}=-\mathbb{1}$, implies

$$
\begin{equation*}
\Gamma_{2}\left(p^{2}\right)=i\left(p^{2}-m_{0}^{2}-M^{2}\left(p^{2}\right)\right) \tag{7.181}
\end{equation*}
$$

In view of (7.178), we identify $\left.\Gamma_{2}\left(p^{2}\right)\right|_{\text {loop }}=-i M^{2}\left(p^{2}\right)$. Together with (7.175) this justifies the above claim that $\Gamma_{2}\left(x_{1}, x_{2}\right)$ consists of the 1- and higher loop 1PI amputated diagrams minus the tree-level propagator.

## 7.9 $\Gamma(\varphi)$ as a quantum effective action and background field method

Diagrams like the ones above illustrate the following interpretation of $\Gamma(\varphi)$ : To compute a connected $n$-point function in the full quantum theory, we use the tree-level Feynman rules but replace each $k$ vertex (with $k \geq 3$ ) of the classical action $S[\phi]$ with the 1 PI amputated $k$-vertex $\Gamma_{k}$ as encoded in $\Gamma[\varphi]$ and replace the free propagator with $G_{2}$, the fully resummed propagator. This justifies the name 1PI effective action of $\Gamma[\varphi]$ :

## Replacing $S[\phi]$ by $\Gamma[\varphi]$ and computing at tree-level gives the full quantum theory.

This interpretation of $\Gamma[\varphi]$ can also be seen in many different ways, one of them being the following: Define the object $\tilde{W}[J, \hbar]$ via ${ }^{8}$

$$
\begin{equation*}
e^{i \tilde{W}[J, \hbar]}:=\int \mathcal{D} \varphi e^{\frac{i}{\hbar}(\Gamma[\varphi)+J \cdot \varphi)} \tag{7.182}
\end{equation*}
$$

[^38]and expand
\[

$$
\begin{equation*}
\tilde{W}[J, \hbar]=\sum_{\text {loops } L} \hbar^{L-1} W^{L}[J] . \tag{7.183}
\end{equation*}
$$

\]

Now consider the formal limit $\hbar \rightarrow 0$, in which the righthand side becomes just $W^{0} / \hbar$. By the stationary phase argument, the right-hand side of (7.182) is given by the value of the integrand at the saddle-point, i.e. by

$$
\begin{equation*}
e^{\frac{i}{\hbar}(\Gamma[\tilde{\varphi}]+J \cdot \tilde{\varphi})} \tag{7.184}
\end{equation*}
$$

with $\tilde{\varphi}$ such that

$$
\begin{equation*}
\left.\frac{\delta}{\delta \varphi}(\Gamma[\varphi]+J \cdot \varphi]\right)\left.\right|_{\varphi=\tilde{\varphi}}=0 \quad \text { i.e. }\left.\frac{\delta \Gamma}{\delta \varphi}\right|_{\tilde{\varphi}}=-J . \tag{7.185}
\end{equation*}
$$

We omit the tilde from now. Thus

$$
\begin{equation*}
W^{0}[J]=\Gamma[\varphi]+J \cdot \varphi \text { with } \frac{\delta \Gamma}{\delta \varphi}=-J, \tag{7.186}
\end{equation*}
$$

where the righthand side is a functional of $J$ only. But this is just the inverse Légendre transformation that defined $\Gamma[\varphi]=W[J]-J \cdot \varphi$ with $\frac{\delta \Gamma}{\delta \varphi}=-J$. Thus $W^{0}[J]=W[J]$. Now, $W[J]$ is the full quantum object, and we thus find that

$$
\begin{equation*}
e^{\frac{i}{\hbar} W[J]}=\frac{1}{Z[0]} \int \mathcal{D} \phi e^{\frac{i}{\hbar}(S[\phi]+J \cdot \phi)}=e^{\frac{i}{\hbar}(\Gamma[\varphi]+J \cdot \varphi)}, \tag{7.187}
\end{equation*}
$$

where $J$ and $\Gamma[\varphi]$ are related via $\frac{\delta \Gamma}{\delta \varphi}=-J$. We can interpret this result as follows:

- The integral $\frac{1}{Z[0]} \int \mathcal{D} \phi$ is responsible for the quantum fluctuations.
- Replacing $S[\phi]$ by $\Gamma[\varphi]$ (with $\varphi$ such that $\frac{\delta \Gamma}{\delta \varphi}=-J$ ) and working at tree-level (i.e. no $\int \mathcal{D} \phi$ ) anymore) gives the full result.

In particular $\Gamma$ and $S$ are the same functionals at tree-level,

$$
\begin{equation*}
\Gamma[\varphi]=S[\varphi]+\hbar K[\varphi] \text { for some } K[\varphi] . \tag{7.188}
\end{equation*}
$$

The equation

$$
\begin{equation*}
\frac{\delta \Gamma}{\delta \varphi}=-J, \quad \text { i.e. } \frac{\delta \Gamma}{\delta \varphi}=0 \quad \text { for } \quad J=0 \tag{7.189}
\end{equation*}
$$

is the quantum effective equation of motion replacing the classical $\frac{\delta S}{\delta \phi}=0$ in the full quantum theory.

## Background field method

Consider again

$$
\begin{equation*}
e^{\frac{i}{\hbar} W[J]}=\int \tilde{\mathcal{D}} \phi e^{\frac{i}{\hbar}(S[\phi]+J \cdot \phi)} \tag{7.190}
\end{equation*}
$$

where $\int \tilde{\mathcal{D}} \phi \equiv \frac{1}{Z[0]} \int \mathcal{D} \phi$ is the normalised measure such that

$$
\begin{equation*}
\left.\int \tilde{\mathcal{D}} \phi e^{\frac{i}{\hbar}(S[\phi]+J \cdot \phi)}\right|_{J=0}=1 . \tag{7.191}
\end{equation*}
$$

If we plug in the relation $W[J]=\Gamma[\varphi]+J \cdot \varphi$ and $\frac{\delta \Gamma}{\delta \varphi}=-J$ we arrive at

$$
\begin{equation*}
e^{\frac{i}{\hbar} \Gamma[\varphi]}=\int \tilde{\mathcal{D}} \phi e^{\frac{i}{\hbar}\left(S[\phi]-\frac{\delta \Gamma}{\delta \varphi}(\phi-\varphi)\right)} . \tag{7.192}
\end{equation*}
$$

Now define the field

$$
\begin{equation*}
f(x):=\phi(x)-\varphi(x), \text { i.e. } \phi(x)=\varphi(x)+f(x) . \tag{7.193}
\end{equation*}
$$

Since $\varphi(x)$ is the vacuum expectation value of the quantum operator $\hat{\phi}(x)$, we view $f(x)$ as the quantum fluctuation around the background $\varphi(x)$. The relation

$$
\begin{equation*}
e^{\frac{i}{\hbar} \Gamma[\varphi]}=\int \tilde{\mathcal{D}} f e^{\frac{i}{\hbar}\left(S[\varphi+f]-\frac{\delta \rho}{\delta \varphi} \cdot f\right)} \tag{7.194}
\end{equation*}
$$

shows that in order to determine $\Gamma[\varphi]$ we must integrate out the quantum fluctuation $f(x)$.
To evaluate this integral we expand

$$
\begin{equation*}
S[\varphi+f]=S[\varphi]+\frac{\delta S}{\delta \varphi} \cdot f+\frac{1}{2} f \cdot \frac{\delta^{2} S}{\delta \varphi^{2}} \cdot f+\text { higher terms } \tag{7.195}
\end{equation*}
$$

and with $\Gamma=S+\hbar \cdot K$ and thus

$$
\begin{equation*}
-\frac{\delta \Gamma}{\delta \varphi} \cdot f=-\frac{\delta S}{\delta \varphi} \cdot f-\hbar \frac{\delta K}{\delta \varphi} \cdot f \tag{7.196}
\end{equation*}
$$

we obtain

$$
\begin{equation*}
e^{\frac{i}{\hbar} \Gamma[\varphi]}=\int \tilde{\mathscr{D}} f e^{\frac{i}{\hbar}\left(S[\varphi]+\frac{1}{2} f \cdot \frac{\delta^{2} \mathcal{S}}{\delta \varphi^{2}} \cdot f-\hbar \frac{\delta \mathcal{L}}{\delta \varphi} \cdot f+O\left(f^{3}\right)\right)} \tag{7.197}
\end{equation*}
$$

Thus altogether

$$
\begin{equation*}
\Gamma[\varphi]=S[\varphi]-i \hbar \ln \left(\int \tilde{\mathcal{D}} f e^{\frac{i}{h}\left(\frac{1}{2} f \cdot \frac{\delta^{2} S}{\delta \varphi^{2}} \cdot f-\hbar \frac{\delta K}{\delta \varphi} \cdot f+O\left(f^{3}\right)\right)}\right) \equiv S[\varphi]+\hbar K[\varphi] . \tag{7.198}
\end{equation*}
$$

Since $K[\varphi]$ appears on both sides, one has to solve this equation perturbatively. In particular the 1-loop correction to $\Gamma$ is due to the lowest fluctuation terms in the integrand of $\int \tilde{\mathcal{D}} f$. Since $K[\varphi]$ is by itself already 1-loop, the relevant piece is only the quadratic one and

$$
\begin{equation*}
K^{(1-\text { loop })}=-i \ln \int \tilde{\mathcal{D}} f e^{-\frac{1}{2} f \cdot\left(-\frac{i}{\hbar} \frac{\delta^{2} s}{\delta \psi^{2}}\right) \cdot f} \tag{7.199}
\end{equation*}
$$

At 2-loop order we then use the 1-loop result for $K[\varphi]$ together with higher-order terms. As an example consider $\phi^{4}$-theory with classical action

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left[-\frac{1}{2} \phi\left(\partial^{2}+m^{2}\right) \phi-\frac{\lambda}{4!} \phi^{4}\right] \tag{7.200}
\end{equation*}
$$

and compute $\Gamma\left[\varphi_{0}\right]$ up to 1-loop order with $\varphi_{0}=$ const. The computation will be done in the tutorials. If we define the effective potential $V\left(\varphi_{0}\right)$ via

$$
\begin{equation*}
\Gamma\left[\phi_{0}\right]=-\mathrm{Vol}_{\mathbb{R}^{1,3}} \cdot V\left(\varphi_{0}\right) \tag{7.201}
\end{equation*}
$$

and expand

$$
\begin{equation*}
V\left(\varphi_{0}\right)=V^{(0)}\left(\varphi_{0}\right)+V^{(1-\text { loop })}\left(\varphi_{0}\right)+\ldots, \tag{7.202}
\end{equation*}
$$

the result is

$$
\begin{equation*}
V^{(0)}=\frac{1}{2} m^{2} \varphi_{0}^{2}+\frac{\lambda}{4!} \varphi_{0}^{4} \tag{7.203}
\end{equation*}
$$

and $V^{(1-\text { loop })}\left(\varphi_{0}\right)$ is given by the Coleman-Weinberg-Potential

$$
\begin{equation*}
V^{(1-\mathrm{loop})}\left(\varphi_{0}\right)=\frac{1}{2} \int \frac{\mathrm{~d}^{4} k}{(2 \pi)^{4} i} \ln \left(\frac{m^{2}+\frac{\lambda}{2} \varphi_{0}^{2}-k^{2}-i \epsilon}{m^{2}-k^{2}-i \epsilon}\right) . \tag{7.204}
\end{equation*}
$$

The denominator is the contribution from the vacuum bubbles of the free theory, which must be removed with our normalisation of the path-integral.

### 7.10 Euclidean QFT and statistical field theory

Let us briefly summarize our findings so far for the path integral in Minkowski space and collect the corresponding expressions for the Euclidean theory. We then point out a beautiful formal analogy between Euclidean Quantum Field Theory in 4 dimensions and Classical Statistical Field Theory in 3 dimensions.
In Lorentzian signature we have derived the following expressions (where we reinstate $\hbar$ ): For the scalar field with classical action

$$
\begin{equation*}
S[\phi]=\int \mathrm{d}^{4} x\left(\frac{1}{2}\left(\partial_{\mu} \phi\right)^{2}-\frac{1}{2} m^{2} \phi^{2}+\mathcal{L}_{\mathrm{int}}\right) \tag{7.205}
\end{equation*}
$$

the generating functional for the correlation functions is defined as

$$
\begin{equation*}
Z[J]=\int \mathcal{D} \phi e^{\frac{i}{\hbar}(S[\phi]+\phi \cdot J)} \tag{7.206}
\end{equation*}
$$

so that

$$
\begin{equation*}
\left\langle T \prod_{i} \phi\left(x_{i}\right)\right\rangle=\left.\frac{1}{Z[0]} \prod_{i} \frac{\hbar \delta}{i \delta J\left(x_{i}\right)} Z[J]\right|_{J=0} . \tag{7.207}
\end{equation*}
$$

Our conventions for the effective action are furthermore

$$
\begin{equation*}
e^{\frac{i}{\hbar} W[J]}=\frac{Z[J]}{Z[0]}, \quad \Gamma[\varphi]=W\left[J_{\varphi}\right]-\varphi \cdot J_{\varphi} . \tag{7.208}
\end{equation*}
$$

As discussed, we can pass from Lorentzian to Euclidean signature via a Wick rotation replacing $x^{0} \rightarrow-i x_{E}^{4}$. In Euclidean space $\mathbb{R}^{4}$ with

$$
\begin{equation*}
\left(x_{E}^{1}, x_{E}^{2}, x_{E}^{3}, x_{E}^{4}\right)=\left(x^{2}, x^{2}, x^{3}, i x^{0}\right) \tag{7.209}
\end{equation*}
$$

the analogous expressions are then

$$
\begin{equation*}
Z_{E}[J]=\int \mathcal{D} \phi e^{-\frac{1}{\hbar}\left(S_{E}[\phi]-\phi \cdot J\right)} \tag{7.210}
\end{equation*}
$$

with

$$
\begin{equation*}
S_{E}[\phi]:=\int \mathrm{d}^{4} x_{E}\left(\frac{1}{2}\left(\partial_{E} \phi\right)^{2}+\frac{1}{2} m^{2} \phi^{2}-\mathcal{L}_{\mathrm{int}}\right) \tag{7.211}
\end{equation*}
$$

and

$$
\begin{equation*}
\phi \cdot J=\int \mathrm{d}^{4} x_{E} \phi\left(x_{E}\right) J\left(x_{E}\right), \quad\left(\partial_{E} \phi\right)^{2}=\sum_{i=1}^{4}\left(\frac{\partial}{\partial x_{E}^{i}} \phi\right)^{2} . \tag{7.212}
\end{equation*}
$$

Furthermore

$$
\begin{equation*}
\left\langle T \prod_{i} \phi\left(x_{E} i\right)\right\rangle=\frac{1}{Z_{E}[0]} \prod_{i} \frac{\hbar \delta}{\delta J\left(x_{E i}\right)} Z_{E}[J] \tag{7.213}
\end{equation*}
$$

and

$$
\begin{equation*}
e^{-\frac{1}{\hbar} W_{E}[J]}:=\frac{Z_{E}[J]}{Z_{E}[0]}, \quad \Gamma_{E}[\varphi]:=W_{E}\left[J_{\varphi}\right]+\varphi \cdot J_{\varphi}, \quad \frac{\delta \Gamma_{E}}{\delta \varphi}=J . \tag{7.214}
\end{equation*}
$$

An executive summary of what Quantum Field Theory is about can be given as follows (working in Euclidean signature):

- We start with a classical field $\phi(x)$ and associate to it the classical action $S_{E}[\phi]$. The field $\phi(x)$ arises as the continuum limit description of a system of $N$ harmonic oscillators.
- In the classical limit, formally $\hbar \rightarrow 0$, the field $\phi(x)$ takes a definite value given by the classical equation of motion

$$
\begin{equation*}
\frac{\delta S_{E}}{\delta \phi}=0 \tag{7.215}
\end{equation*}
$$

- For $\hbar$ finite, quantum fluctuations arise. These are encoded in $Z_{E}[J]$, where the path integral takes into account all quantum fluctuations of the field.
- What we can compute in the quantum theory are quantum expectation values or correlation functions. In particular the quantum expectation value of the field in the presence of a source $J$ is

$$
\begin{equation*}
\varphi(x)_{J}:=\langle\phi(x)\rangle_{J}=\langle\Omega| \hat{\phi}(x)|\Omega\rangle_{J}=\frac{1}{Z[0]} \int \mathcal{D} \phi \phi(x) e^{-\frac{1}{\hbar}\left(S_{E}[\phi]-\phi \cdot J\right)} . \tag{7.216}
\end{equation*}
$$

- With our definition of $W_{E}[J]$ we can compute this as

$$
\begin{equation*}
\varphi(x)_{J}=-\frac{\delta W_{E}}{\delta J(x)} \tag{7.217}
\end{equation*}
$$

- In terms of the Légendre transform $\Gamma[\varphi]$ we have

$$
\begin{equation*}
\frac{\delta \Gamma_{E}}{\delta \varphi}=J \tag{7.218}
\end{equation*}
$$

$\Gamma_{E}$ gives a quantum effective action after integrating out the quantum fluctuations in the sense that

$$
\begin{equation*}
e^{-\frac{1}{\hbar} \Gamma_{E}[\varphi]}=\int \tilde{\mathcal{D}} f e^{-\frac{1}{\hbar}\left(S_{E}[\varphi+f]-\frac{\delta \Gamma_{E}}{\delta_{\varphi}} \cdot f\right)} \tag{7.219}
\end{equation*}
$$

Remarkably, this is formally exactly the same structure as in classical statistical mechanics or more precisely in its continuum limit, classical statistical field theory, in 3 spatial dimensions:

- Start with a classical field describing some classical observable, e.g. with $s(\vec{x})$ defined as the magnetisation density of a 3D ferromagnet with energy

$$
\begin{equation*}
\int \mathrm{d}^{3} x \mathcal{H}(\vec{x})=\int \mathrm{d}^{3} x \frac{1}{2}(\nabla s(\vec{x}))^{2} . \tag{7.220}
\end{equation*}
$$

We can think of this as the continuum limit of a system of spins $s_{n}$ with nearest neighbour interactions.

- At temperature $T=0$, the system is in its ground state given by minimizing the energy (in this case given by $s(\vec{x})$ constant).
- For finite $T$ there are thermal fluctuations of $s(\vec{x})$. These are encoded in the partition function

$$
\begin{equation*}
Z[H]=\int \mathcal{D} s(\vec{x}) e^{-\beta \int \mathrm{d}^{3} x(\mathcal{H}(s)-H(\vec{x}) s(\vec{x}))}, \tag{7.221}
\end{equation*}
$$

where $\beta=\left(k_{B} T\right)^{-1}$ and $H(\vec{x})$ is an external source coupling to $s(\vec{x})$ (in the present case of a ferromagnet this would be a magnetic field).

- The classical expectation value of $s(\vec{x})$ for an ensemble of systems is

$$
\begin{equation*}
\langle s(\vec{x})\rangle_{H}=\frac{1}{Z[H]} \int \mathcal{D} s(\vec{x}) s(\vec{x}) e^{-\beta \int \mathrm{d}^{3} x(\mathcal{H}-H s)} . \tag{7.222}
\end{equation*}
$$

- Define the free energy - the analogue of $W_{E}[J]$ - as

$$
\begin{equation*}
F[H]=-\frac{1}{\beta} \ln Z, \tag{7.223}
\end{equation*}
$$

i.e.

$$
\begin{equation*}
e^{-\beta F[H]}=\int \mathcal{D} s(\vec{x}) e^{-\beta \int \mathrm{d}^{3} x(\mathcal{H}-H s)} \tag{7.224}
\end{equation*}
$$

Then

$$
\begin{equation*}
M(\vec{x}):=\langle s(\vec{x})\rangle_{H}=-\frac{\delta F}{\delta H(\vec{x})} . \tag{7.225}
\end{equation*}
$$

- The Gibbs free energy is the Légendre transform of $F[H]$,

$$
\begin{equation*}
G[M]=F\left[H_{M}\right]+M \cdot H_{M}, \tag{7.226}
\end{equation*}
$$

viewed as a functional of $M$. Then the average magnetisation is determined via the equation

$$
\begin{equation*}
\frac{\delta G}{\delta M}=H . \tag{7.227}
\end{equation*}
$$

Due to this formal analogy, many problems in statistical field theory are solved in exactly the same manner as in Euclidean QFT. In particular the structure of perturbation theory and important concepts concerning symmetry breaking and renormalisation agree.
This can be summarized as follows:

| 4D QFT | 3D statistical theory |
| :---: | :---: |
| $\phi\left(x_{E}\right)$ | $s(\vec{x})$ |
| quantum fluctuation | thermal fluctuation |
| $\hbar$ | $k_{B} T$ |
| $e^{-\frac{1}{\hbar} W_{E}[J]}$ | $e^{-\beta F[H]}$ |
| $\varphi_{J}=-\frac{\delta W_{E}}{\delta J(\vec{x})}$ | $M=-\frac{\delta F}{\delta H(\vec{x})}$ |
| $\frac{\delta \Gamma_{E}}{\delta \varphi}=J$ | $\frac{\delta G}{\delta M}=H$ |

### 7.11 Grassman algebra calculus

We would now like to apply the path integral quantization method to systems with fermions. Consider therefore a Dirac fermion field operator $\hat{\psi}^{A}(t, \vec{x})$ with spin indices $A=1,2,3,4$ in the Heisenberg picture and recall the anti-commutation relations

$$
\begin{equation*}
\left\{\hat{\psi}^{A}(t, \vec{x}), \hat{\psi}_{B}^{\dagger}(t, \vec{y})\right\}=\delta_{B}^{A} \delta^{(3)}(\vec{x}-\vec{y}) \tag{7.228}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\{\hat{\psi}^{A}(t, \vec{x}), \hat{\psi}^{B}(t, \vec{y})\right\}=0=\left\{\hat{\psi}_{A}^{\dagger}(t, \vec{x}), \hat{\psi}_{B}^{\dagger}(t, \vec{y})\right\} . \tag{7.229}
\end{equation*}
$$

We seek a path integral representation for the transition function $\left\langle\psi\left(\vec{x}_{F}\right), t_{f} \mid \psi\left(\vec{x}_{I}\right), t_{I}\right\rangle$ where the state $|\psi(\vec{x}), t\rangle$ is defined via

$$
\begin{equation*}
\hat{\psi}^{A}(t, \vec{x})|\psi(\vec{x}), t\rangle=\psi^{A}(t, \vec{x})|\psi(\vec{x}), t\rangle . \tag{7.230}
\end{equation*}
$$

Here the non-hatted expression $\psi^{A}(t, \vec{x})$ represents a classical field. We face the immediate problem that consistency with (7.229) requires that the classical field $\psi^{A}(t, \vec{x})$ satisfy

$$
\begin{equation*}
\psi^{A}(t, \vec{x}) \psi^{B}(t, \vec{y})=-\psi^{B}(t, \vec{y}) \psi^{A}(t, \vec{x}) . \tag{7.231}
\end{equation*}
$$

Thus $\psi^{A}(t, \vec{x})$ cannot take values in $\mathbb{C}$ - rather we need the notion of anti-commuting numbers.

Consider first the case of a finite number of degrees of freedom by putting the theory, as usual, on a lattice of finite lattice spacing corresponding to a transition

$$
\begin{equation*}
\psi^{A}(t, \vec{x}) \rightarrow \psi_{i}^{A}(t) \tag{7.232}
\end{equation*}
$$

and let us omit the spinor indices for the time being. Then we search for classical objects $\psi_{i}(t)$ with the property

$$
\begin{equation*}
\psi_{i}(t) \psi_{j}(t)=-\psi_{j}(t) \psi_{i}(t) \tag{7.233}
\end{equation*}
$$

These $\psi_{i}$ take values in a so-called Grassmann-algebra $\mathbb{A}$ defined as follows:

- Let $\theta_{i}, i=1, \ldots, n$, be a basis of an $n$-dimensional complex vector space $V$, i.e. we have the notion of scalar multiplication $a \theta_{i}=\theta_{i} a, a \in \mathbb{C}$ and vector addition $a \theta_{i}+b \theta_{j} \in V$.
- Define then a bilinear anti-commutative multiplication

$$
\begin{equation*}
\theta_{i} \theta_{j}=-\theta_{j} \theta_{i} \tag{7.234}
\end{equation*}
$$

as a map from

$$
\begin{equation*}
V \times V \mapsto \Lambda^{2} V, \tag{7.235}
\end{equation*}
$$

where $\Lambda^{2} V$ is the antisymmetric tensor product of $V$.

- Iteratively we can build higher rank anti-symmetric product spaces $\Lambda^{k} V$ up to $\Lambda^{n} V$.

Then the Grassmann (or exterior) algebra is defined as the space

$$
\begin{equation*}
\mathbb{A}=\bigoplus_{k=0}^{n} \Lambda^{k} V=\mathbb{C} \oplus V \oplus \Lambda^{2} V \oplus \ldots \oplus \Lambda^{n} V \tag{7.236}
\end{equation*}
$$

A typical element of $\mathbb{A}$ is of the form

$$
\begin{equation*}
a+a_{i} \theta_{i}+\frac{1}{2} a_{i j} \theta_{i} \theta_{j}+\ldots+\frac{1}{n!} a_{i_{1} \ldots i_{n}} \theta_{i_{1} \ldots \theta_{i_{n}}} \tag{7.237}
\end{equation*}
$$

with completely antisymmetrc coefficients $a_{i j}=-a_{j i}$ etc. By abuse of notation, the $\theta_{i}$ are called Grassmann numbers - though calling them 'numbers' makes things look more exotic than they are as the $\theta_{i}$ should really be viewed simply as elements of the abstract Grassmann algebra $\mathbb{A}$. In particular they are nilpotent,

$$
\begin{equation*}
\theta_{i}^{2}=0 \tag{7.238}
\end{equation*}
$$

An important example of a Grassmann algebra is the space $\Omega$ of differential forms defined on an $n$-dimensional manifold endowed with the structure wedge (or exterior product).
A Grassmann algebra is a graded algebra:
An element of $\Lambda^{2 k} V$ has grade (or degree) $s=0$ (it is said to be even), an element of $\Lambda^{2 k+1} V$ has grade (or degree) $s=1$ (it is said be odd).

Note that a general element of $\mathbb{A}$ can always be written as the sum of an even and an odd element. Two elements $A, B \in \mathbb{A}$ of definite degree $s_{A}$ and $s_{B}$ satisfy

$$
\begin{equation*}
A B=(-1)^{s_{A} s_{B}} B A . \tag{7.239}
\end{equation*}
$$

Thus even elements are commuting and therefore sometimes called bosonic, whereas odd elements are dubbed fermionic due to their anti-commuting nature.

## Differentiation and integration

To set up a notion of calculus on the space of functions

$$
\begin{align*}
f: \mathbb{A} & \rightarrow \mathbb{C},  \tag{7.240}\\
\underline{\theta} & \rightarrow f(\underline{\theta})=a+a_{i} \theta_{i}+\ldots+\frac{1}{n!} a_{i_{1} \ldots i_{n}} \theta_{i_{1 i} i_{n}}, \tag{7.241}
\end{align*}
$$

one defines differentiation with respect to $\theta_{i}$ as follows:

- $\frac{\partial}{\partial \theta_{i}} \theta_{j}=\delta_{i j}, \quad \frac{\partial}{\partial \theta_{j}} a=0 \quad \forall a \in \mathbb{C} ;$
- $\frac{\partial}{\partial \theta_{i}}\left(a_{1} f_{1}(\underline{\theta})+a_{2} f_{2}(\underline{\theta})\right)=a_{1} \frac{\partial}{\partial \theta_{1}} f_{1}(\underline{\theta})+a_{2} \frac{\partial}{\partial \theta_{2}} f_{2}(\underline{\theta})$ (linearity);
- If $f_{1}(\underline{\theta})$ has definite grade $s$, then the graded Leibniz rule holds,

$$
\begin{equation*}
\frac{\partial}{\partial \theta_{i}}\left(f_{1}(\underline{\theta}) f_{2}(\underline{\theta})\right)=\left(\frac{\partial}{\partial \theta_{i}} f_{1}(\underline{\theta})\right) f_{2}(\underline{\theta})+(-1)^{s} f_{1}(\underline{\theta}) \frac{\partial}{\partial \theta_{i}} f_{2}(\underline{\theta}) . \tag{7.242}
\end{equation*}
$$

For general $f_{1}(\underline{\theta})$ we decompose $f_{1}(\underline{\theta})$ into its even and odd part and then apply the graded Leibniz rule by linearity.

For example

$$
\begin{equation*}
\frac{\partial}{\partial \theta_{i}}\left(\theta_{j} \theta_{k}\right)=\delta_{i j} \theta_{k}-\delta_{i k} \theta_{j} . \tag{7.243}
\end{equation*}
$$

Next we define an integral $I[f(\underline{\theta})]$ as a functional of $f(\underline{\theta})$. Our applications to the path integral will make it evident that it is useful to define this functional such that it shares two key properties of the bosonic integral (in one variable for simplicity)

$$
\begin{equation*}
I[f(x)]=\int_{-\infty}^{\infty} f(x) \mathrm{d} x, \tag{7.244}
\end{equation*}
$$

namely

- $I\left[a_{1} f_{1}(x)+a_{2} f_{2}(x)\right]=a_{1} I\left[f_{1}(x)\right]+a_{2} I\left[f_{2}(x)\right]$ and
- $I[f(x+a)]=I[f(x)]$, i.e. translation invariance in the integration variable.

Consider first a Grassmann algebra with $n=1$ such that

$$
\begin{equation*}
I\left[f(\theta]=\int \mathrm{d} \theta f(\theta)=\int \mathrm{d} \theta(a+b \theta) .\right. \tag{7.245}
\end{equation*}
$$

In this case translation invariance just means invariance under a shift of the integration variable $\theta \rightarrow$ $\theta+c$ and implies, together with linearity,

$$
\begin{equation*}
\int \mathrm{d} \theta c=0 \quad \forall c \in \mathbb{C} . \tag{7.246}
\end{equation*}
$$

We normalise $\int \mathrm{d} \theta \theta=1$ such that

$$
\begin{equation*}
\int \mathrm{d} \theta[a+b \theta]=b=\frac{\partial}{\partial \theta}(a+b \theta) . \tag{7.247}
\end{equation*}
$$

For a Grassmann algebra based on Grassmann numbers $\theta_{i}, i=1, \ldots, n$, generalized translation invariance implies

$$
\begin{equation*}
\int \mathrm{d} \theta_{i} \theta_{j}=\delta_{i j}=\frac{\partial}{\partial \theta_{i}} \theta_{j} . \tag{7.248}
\end{equation*}
$$

Then the Grassmann measure is defined as $\mathrm{d}^{n} \theta:=\mathrm{d} \theta_{n} \mathrm{~d} \theta_{n-1} \ldots \mathrm{~d} \theta_{1}$. Because of

$$
\begin{equation*}
\mathrm{d} \theta_{i} \mathrm{~d} \theta_{j}=-\mathrm{d} \theta_{j} \mathrm{~d} \theta_{i} \tag{7.249}
\end{equation*}
$$

we have to be aware of the order in $\mathrm{d}^{n} \theta$. Thus $\int \mathrm{d}^{n} \theta \theta_{1} \theta_{2} \ldots \theta_{n}=1$ and more generally

$$
\begin{equation*}
\int \mathrm{d}^{n} \theta \theta_{i_{1}} \ldots \theta_{i_{n}}=\epsilon_{i_{1} \ldots i_{n}} \tag{7.250}
\end{equation*}
$$

with $\epsilon_{i_{1} \ldots i_{n}}$ the totally anti-symmetric $\epsilon$-tensor normalised as $\epsilon_{1 \ldots, n}=1$. In particular for $f(\underline{\theta})=$ $a+a_{i} \theta_{i}+\ldots+\frac{1}{n!} a_{i_{1} \ldots i_{n}} \theta_{i_{1} \ldots i_{n}}$ we find

$$
\begin{equation*}
\int \mathrm{d}^{n} \theta f(\underline{\theta})=\frac{1}{n!} a_{i_{1} \ldots i_{n}} \epsilon_{i_{1} \ldots i_{n}}=a_{123 \ldots n} \tag{7.251}
\end{equation*}
$$

Only those terms contribute for which the Grassmann integral is 'saturated' such that each $d \theta_{i}$ is paired with precisely one $\theta_{i}$.
Note that the above definitions ensure that upon integration by parts no boundary terms are picked up,

$$
\begin{equation*}
\int \mathrm{d}^{n} \theta \frac{\partial}{\partial \theta_{i}} f(\underline{\theta})=0 . \tag{7.252}
\end{equation*}
$$

As an application we compute the Jacobian that arises for a linear change of Grassmann integration variables. To begin with, if $A$ is an $(n \times n)$-matrix, then we can define $\underline{\theta}^{\prime}=A \underline{\theta}$ and find

$$
\begin{align*}
\int \mathrm{d}^{n} \theta f\left(\underline{\theta}^{\prime}\right)=\int \mathrm{d}^{n} \theta f(A \underline{\theta}) & =\int \mathrm{d}^{n} \theta a_{12 \ldots n} A_{1 i_{1}} \ldots A_{n i_{n}} \theta_{i_{1}} \ldots \theta_{i_{n}}  \tag{7.253}\\
& =a_{12 \ldots n} \underbrace{A_{1 i_{1} \ldots . A_{n i_{n}} \epsilon_{11} \ldots i_{n}}}_{=\operatorname{det} A}=\operatorname{det} \mathrm{A} \int \mathrm{~d}^{n} \theta f(\underline{\theta}) . \tag{7.254}
\end{align*}
$$

In the first line we used (7.251). Thus

$$
\begin{equation*}
\int \mathrm{d}^{n} \theta f\left(\underline{\theta}^{\prime}\right)=\operatorname{det} A \int \mathrm{~d}^{n} \theta f(\underline{\theta})=\operatorname{det} A \int \mathrm{~d}^{n} \underline{\theta}^{\prime} f\left(\underline{\theta}^{\prime}\right) \tag{7.255}
\end{equation*}
$$

and consequently $\mathrm{d}^{n} \theta=\operatorname{det} A \mathrm{~d}^{n} \theta^{\prime}$, i.e.

$$
\begin{equation*}
\underline{\theta}^{\prime}=A \underline{\theta} \Rightarrow \mathrm{~d}^{n} \underline{\theta}^{\prime}=(\operatorname{det} A)^{-1} \mathrm{~d}^{n} \underline{\theta} . \tag{7.256}
\end{equation*}
$$

Compare this with the bosonic case $\mathrm{d}^{n}(A \cdot x)=\|\operatorname{det} A\| \mathrm{d}^{n} x$.

## Application: Fermionic Gauss integral

Consider $\theta_{i}, i=1, \ldots, n$, with $n=2 m$ and let us compute

$$
\begin{equation*}
I=\int \mathrm{d}^{n} \theta e^{\frac{1}{2} \theta_{i} A_{i j} \theta_{j}}, \quad A_{i j}=-A_{j i} . \tag{7.257}
\end{equation*}
$$

To saturate the Grassmann measure we pick up the term to $m$-th order in the expansion of the exponential and find

$$
\begin{align*}
I & =\frac{1}{2^{m} m!} \int \mathrm{d}^{n} \theta A_{i_{1} i_{2}} \ldots A_{i_{n-1} i_{n}} \theta_{i_{1}} \ldots \theta_{i_{n-1}} \theta_{i_{n}}  \tag{7.258}\\
& =\frac{1}{2^{m} m!} \epsilon_{i_{1} \ldots i_{n}} A_{i_{1} i_{2}} \ldots A_{i_{n-1} i_{n}}:=\operatorname{Pf}(A),
\end{align*}
$$

where $\operatorname{Pf}(A)$ is the so-called Pfaffian. Note the difference in the definition of the Pfaffian and of the determinant,

$$
\begin{equation*}
\operatorname{det} A=\frac{1}{n!} \epsilon_{i_{1} \ldots i_{n}} \epsilon_{j_{1} \ldots j_{n}} A_{i_{1} j_{1}} \ldots A_{i_{n} j_{n}} . \tag{7.259}
\end{equation*}
$$

One can show with standard methods of linear algebra that $(\operatorname{Pf}(A))^{2}=\operatorname{det} A$ and take the positive sign,

$$
\begin{equation*}
\operatorname{Pf}(A)=\sqrt{\operatorname{det} A} \tag{7.260}
\end{equation*}
$$

## Complexification

So far the coefficients $a \in \mathbb{C}$, but $\theta_{i}$ were treated as real, or equivalently as one anti-commuting degree of freedom (which is the fermionic counterpart of 1 bosonic real number). One can also consider complex $\theta_{i}$ and $\theta_{i}^{*}$, viewed as independent degrees of freedom. E.g. one can build such complex Grassmann variables $\theta_{i}, i=1, \ldots, n$, out of real Grassmann variables $\omega_{j}, j=1, \ldots, 2 n$ via

$$
\begin{equation*}
\theta_{i}=\omega_{2 i-1}+i \omega_{2 i}, \quad \theta_{i}^{*}=\omega_{2 i-1}-i \omega_{2 i} . \tag{7.261}
\end{equation*}
$$

Alternatively, we simply declare $\theta_{i}$ and $\theta_{i}^{*}$ to denote independent objects and define the complex conjugate of a product via

$$
\begin{equation*}
\left(\theta_{i} \theta_{j}\right)^{*}=\theta_{j}^{*} \theta_{i}^{*} . \tag{7.262}
\end{equation*}
$$

We furthermore define the measure

$$
\begin{equation*}
\mathrm{d}^{n} \theta \mathrm{~d}^{n} \theta^{*}:=\prod_{i=1}^{n} \mathrm{~d} \theta_{i} \mathrm{~d} \theta_{i}^{*}=\mathrm{d} \theta_{1} \mathrm{~d} \theta_{1}^{*} \ldots \mathrm{~d} \theta_{n} \mathrm{~d} \theta_{n}^{*}=\mathrm{d} \theta_{n} \mathrm{~d} \theta_{n}^{*} \ldots \mathrm{~d} \theta_{1} \mathrm{~d} \theta_{1}^{*} \tag{7.263}
\end{equation*}
$$

The complex Gaussian integral then reads

$$
\begin{align*}
\int \mathrm{d} \theta \mathrm{~d} \theta^{*} e^{\theta_{i}^{*} B_{i j} \theta_{j}} & =\frac{1}{n!} \int \mathrm{d} \theta_{n} \mathrm{~d} \theta_{n}^{*} \ldots \mathrm{~d} \theta_{1} \mathrm{~d} \theta_{1}^{*} B_{i_{1} j_{1} \ldots B_{i_{n} j_{n}}} \theta_{i_{1}}^{*} \theta_{j_{1}} \ldots \theta_{i_{n}}^{*} \theta_{j_{n}}  \tag{7.264}\\
& =\frac{1}{n!} \epsilon_{i_{1} \ldots i_{n}} \epsilon_{j_{1} \ldots j_{n}} B_{i_{1} j_{1}} \ldots B_{i_{n} j_{n}}=\operatorname{det} B .
\end{align*}
$$

Let us compare our results for bosonic and fermionic Gaussian integrals:

- Real case:

$$
\begin{align*}
\int \mathrm{d}^{n} x e^{-\frac{1}{2} x_{i} A_{i j} x_{j}} & =\sqrt{\frac{(2 \pi)^{n}}{\operatorname{det} A}} \quad \text { (bosonic integration), }  \tag{7.265}\\
\int \mathrm{d}^{n} \theta e^{\frac{1}{\theta_{i} A_{i j} \theta_{j}}} & =\sqrt{\operatorname{det} A} \quad \text { (fermionic integration). }
\end{align*}
$$

## - Complex case:

$$
\begin{align*}
& \int \mathrm{d}^{n} z \mathrm{~d}^{n} z^{*} e^{-z_{i}^{*} B_{i j} z_{j}}=\frac{(2 \pi)^{n}}{\operatorname{det} B} \quad \text { (bosonic), }  \tag{7.266}\\
& \int \mathrm{d}^{n} \theta^{*} \mathrm{~d}^{n} \theta e^{-\theta_{i}^{*} B_{i j} \theta_{j}}=\operatorname{det} B \quad \text { (fermionic). }
\end{align*}
$$

Note the change in the order of $\mathrm{d}^{n} \theta^{*}$ and $\mathrm{d}^{n} \theta$ in the last term in order to accomodate for the $(-1)$ sign in the exponent.

### 7.12 The fermionic path integral

We are now ready to come back to physics and to address the transition amplitude between two fermionic states. We first consider the quantum mechanics of a single fermion by making the transition

$$
\begin{equation*}
\hat{\psi}^{A}(t, \vec{x}) \rightarrow \hat{\psi}_{i}(t) \tag{7.267}
\end{equation*}
$$

ignoring the spinor components $A$ and taking $i=1$. We then seek the states $|\psi, t\rangle$ such that

$$
\begin{equation*}
\hat{\psi}(t)|\psi, t\rangle=\psi(t)|\psi, t\rangle \tag{7.268}
\end{equation*}
$$

with $\psi(t)$ a complex Grassmann number. Working at fixed time $t$ and suppressing the $t$-dependence we must find the Hilbert space acted upon by $\hat{\psi}$ and $\hat{\psi}^{\dagger}$ with anti-commutation relations

$$
\begin{equation*}
\{\hat{\psi}, \hat{\psi}\}=0=\left\{\hat{\psi}^{\dagger}, \hat{\psi}^{\dagger}\right\}, \quad\left\{\hat{\psi}, \hat{\psi}^{\dagger}\right\}=1 . \tag{7.269}
\end{equation*}
$$

In analogy with the analogous bosonic relations $\left[a, a^{\dagger}\right]=1,[a, a]=0=\left[a^{\dagger}, a^{\dagger}\right]$ this system is called a fermionic harmonic oscillator. We assume (or alternatively show by exploiting boundedness of the Hamiltonian) the existence of a ground state

$$
\begin{equation*}
|0\rangle \text { such that } \hat{\psi}|0\rangle=0 \tag{7.270}
\end{equation*}
$$

and define the state $|1\rangle$ as $\hat{\psi}^{\dagger}|0\rangle=|1\rangle$. By means of the anti-commutation relations

$$
\begin{align*}
\hat{\psi}^{\dagger}|1\rangle & =\hat{\psi}^{\dagger} \hat{\psi}^{\dagger}|0\rangle=0,  \tag{7.271}\\
\hat{\psi}|1\rangle & =\hat{\psi} \hat{\psi}^{\dagger}|0\rangle=\left(1-\hat{\psi}^{\dagger} \hat{\psi}\right)|0\rangle=|0\rangle .
\end{align*}
$$

We now claim that the sought-after eigenstate $|\psi\rangle$ with $\hat{\psi}|\psi\rangle=\psi|\psi\rangle$ takes the form

$$
\begin{equation*}
|\psi\rangle=|0\rangle-\psi|1\rangle . \tag{7.272}
\end{equation*}
$$

To avoid confusion note that the unhatted $\psi$ is a complex Grassmann number, while the hatted $\hat{\psi}, \hat{\psi}^{\dagger}$ are annihilation and creation operators.
Indeed,

$$
\begin{equation*}
\hat{\psi}(|0\rangle-\psi|1\rangle)=0-\hat{\psi} \psi|1\rangle \tag{7.273}
\end{equation*}
$$

Since $|1\rangle=\psi^{\dagger}|0\rangle$ is a Grassmann odd object, we have $\psi|1\rangle=-|1\rangle \psi$. Equivalently, this follows from the relation $\hat{\psi}^{\dagger} \psi=-\psi \hat{\psi}^{\dagger}$. Either way we find

$$
\begin{equation*}
\hat{\psi}|\psi\rangle=\psi \hat{\psi}|1\rangle=\psi|0\rangle=\psi|0\rangle-\psi^{2}|1\rangle=\psi(|0\rangle-\psi|1\rangle)=\psi|\psi\rangle \tag{7.274}
\end{equation*}
$$

We can rewrite the eigenstate as

$$
\begin{equation*}
e^{-\psi \hat{\psi}^{\dagger}}|0\rangle=\left(1-\psi \hat{\psi}^{\dagger}\right)|0\rangle=|0\rangle-\psi|1\rangle \tag{7.275}
\end{equation*}
$$

Thus the eigenstate of the fermionic annihilation operator

$$
\begin{equation*}
\hat{\psi}|\psi\rangle=\psi|\psi\rangle \tag{7.276}
\end{equation*}
$$

is given by

$$
\begin{equation*}
|\psi\rangle=|0\rangle-\psi|1\rangle=e^{-\psi \psi^{\dagger}}|0\rangle \tag{7.277}
\end{equation*}
$$

The state $|\psi\rangle$ is the fermionic analogue of a bosonic coherent state $|\mu\rangle$ defined as the eigenstate of the annihilator,

$$
\begin{equation*}
a|\mu\rangle=\mu|\mu\rangle \quad \text { with } \quad|\mu\rangle=e^{\mu a^{\dagger}}|0\rangle \tag{7.278}
\end{equation*}
$$

The complex conjugate state $\langle\psi|$ defined via $\langle\psi| \hat{\psi}^{\dagger}=\langle\psi| \psi^{*}$ is simply

$$
\begin{equation*}
\langle\psi|=\langle 0|-\langle 1| \psi^{*}=\langle 0| e^{-\hat{\psi} \psi^{*}} \tag{7.279}
\end{equation*}
$$

We will also need the inner product of two fermionic coherent states

$$
\begin{equation*}
\left\langle\psi \mid \psi^{\prime}\right\rangle=\langle 0 \mid 0\rangle+\langle 1| \psi^{*} \psi^{\prime}|1\rangle=1+\psi^{*} \psi^{\prime} \tag{7.280}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\left\langle\psi \mid \psi^{\prime}\right\rangle=e^{\psi^{*} \psi^{\prime} .} \tag{7.281}
\end{equation*}
$$

Finallly the completeness relation in terms of the coherent states is given by

$$
\begin{equation*}
\int \mathrm{d} \psi^{*} \mathrm{~d} \psi|\psi\rangle e^{-\psi^{*} \psi}\langle\psi|=\mathbb{1} \tag{7.282}
\end{equation*}
$$

since the left-hand side is

$$
\begin{align*}
& \int \mathrm{d} \psi^{*} \mathrm{~d} \psi(|0\rangle-\psi|1\rangle)\left(1-\psi^{*} \psi\right)\left(\langle 0|-\langle 1| \psi^{*}\right) \\
& \quad=\int \mathrm{d} \psi^{*} \mathrm{~d} \psi \underbrace{\left(-|0\rangle \psi^{*} \psi\langle 0|+\psi|1\rangle\langle 1| \psi^{*}\right)}_{\left.-\psi^{*} \psi|0\rangle\langle 0|+\psi \psi^{*}|1\rangle\langle 1|\right)}  \tag{7.283}\\
& \quad=|0\rangle\langle 0|+|1\rangle\langle 1|=\mathbb{1}
\end{align*}
$$

where we used that only terms $\sim \psi^{*} \psi$ saturate the Grassmann measure.

Let us now compute the quantum mechanical transition amplitude

$$
\begin{equation*}
\left\langle\psi_{F}, t_{F} \mid \psi_{I}, t_{I}\right\rangle=\left\langle\psi_{F}\right| e^{-i \hat{H}\left(t_{f}-t_{I}\right)}\left|\psi_{I}\right\rangle \tag{7.284}
\end{equation*}
$$

between two fermionic coherent states. We assume that the structure of the Hamiltonian is

$$
\begin{equation*}
\hat{H}=\hat{\psi}^{\dagger} M \hat{\psi} \equiv H\left(\hat{\psi}^{\dagger}, \hat{\psi}\right) \tag{7.285}
\end{equation*}
$$

for some bosonic, i.e. Grassmann even, $M$. In fact this will be sufficient to treat the field theoretic Dirac fermion system we are interested in. As in the derivation of the bosonic path integral we partition the time evolution operator according to

$$
\begin{equation*}
e^{-i \hat{H}\left(t_{F}-t_{I}\right)}=\lim _{N \rightarrow \infty}\left(e^{-i H \delta t}\right)^{N}, \quad \delta t=\frac{t_{F}-t_{I}}{N} \tag{7.286}
\end{equation*}
$$

and insert suitable factors of the identity, this time given by

$$
\begin{equation*}
\mathbb{1}_{j}=\int \mathrm{d} \psi_{j}^{*} \mathrm{~d} \psi_{j}\left|\psi_{j}\right\rangle e^{-\psi_{j}^{*} \psi_{j}}\left\langle\psi_{j}\right| \tag{7.287}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\left\langle\psi_{F}, t_{F} \mid \psi_{I}, t_{I}\right\rangle=\lim _{N \rightarrow \infty} \int \prod_{j=1}^{N} \mathrm{~d} \psi_{j}^{*} \mathrm{~d} \psi_{j}\left\langle\psi_{N+1}\right| e^{-i \hat{H} \delta t}\left|\psi_{N}\right\rangle e^{-\psi_{N}^{*} \psi_{N}}\left\langle\psi_{N}\right| \ldots\left|\psi_{1}\right\rangle e^{-\psi_{1}^{*} \psi_{1}}\left\langle\psi_{1}\right| e^{-i \hat{H} \delta t}\left|\psi_{0}\right\rangle \tag{7.288}
\end{equation*}
$$

with $\left|\psi_{0}\right\rangle \equiv\left|\psi_{I}\right\rangle$ and $\left|\psi_{N+1}\right\rangle \equiv\left|\psi_{F}\right\rangle$. Now

$$
\begin{align*}
\left\langle\psi_{j+1}\right| e^{-i \hat{H} \delta t}\left|\psi_{j}\right\rangle & =\left\langle\psi_{j+1}\right| e^{-i \hat{\psi}^{\dagger} M \hat{\psi} \delta t}\left|\psi_{j}\right\rangle \\
& =\left\langle\psi_{j+1}\right|\left(1-i \hat{\psi}^{\dagger} M \hat{\psi} \delta t\right)\left|\psi_{j}\right\rangle \\
& =\left\langle\psi_{j+1} \mid \psi_{j}\right\rangle-i\left\langle\psi_{j+1}\right| \hat{\psi}^{\dagger} M \hat{\psi} \delta t\left|\psi_{j}\right\rangle \\
& =\left(1-i \psi_{j+1}^{*} M \psi_{j} \delta t\right)\left\langle\psi_{j+1} \mid \psi_{j}\right\rangle  \tag{7.289}\\
& =e^{\psi_{j+1}^{*} \psi_{j}} e^{-i \psi_{j+1}^{*} M \psi_{j} \delta t} \\
& =e^{\psi_{j+1}^{*} \psi_{j+1}} e^{\psi_{j+1}^{*}\left(\psi_{j}-\psi_{j+1}\right)} e^{-i \psi_{j+1}^{*} M \psi_{j} \delta t}
\end{align*}
$$

The factor $e^{\psi_{j+1}^{*} \psi_{j+1}}$ cancels with a corresponding factor $e^{-\psi_{j+1}^{*} \psi_{j+1}}$ from the insertion of the identity, except for one remaining factor $e^{\psi_{N+1}^{*} \psi_{N+1}}$. Thus

$$
\begin{equation*}
\left\langle\psi_{F}, t_{F} \mid \psi_{I}, t_{I}\right\rangle=\lim _{N \rightarrow \infty} \int \prod_{j=1}^{N} \mathrm{~d} \psi_{j}^{*} \mathrm{~d} \psi_{j} e^{\psi_{N+1}^{*} \psi_{N+1}} e^{\sum_{j=0}^{N}\left(-\psi_{j+1}^{*}\left(\psi_{j+1}-\psi_{j}\right)-i \delta t H\left(\psi_{j+1}^{*}, \psi_{j}\right)\right)} \tag{7.290}
\end{equation*}
$$

As $N \rightarrow \infty$ the exponential becomes

$$
\begin{equation*}
e^{i \int \mathrm{~d} t\left(\psi^{*} i \frac{\partial}{\partial t} \psi-H\right)} \tag{7.291}
\end{equation*}
$$

and $e^{\psi_{N+1}^{*} \psi_{N+1}}$ can be absorbed into the measure to arrive at the fermionic path integral for the Quantum Mechanics of a single complex fermion,

$$
\begin{equation*}
\left\langle\psi_{F}, t_{F} \mid \psi_{I}, t_{I}\right\rangle=\int \mathcal{D} \psi^{*} \mathcal{D} \psi e^{i \int_{t_{I}}^{t_{F}} \mathrm{~d} t\left(\psi^{*} i \frac{\partial \psi}{\partial t}-H\right)} \tag{7.292}
\end{equation*}
$$

To generalise this result to the Quantum Field Theory of a Dirac fermion we consider the classical action

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x \bar{\psi}(x)\left(i \gamma^{\mu} \partial_{\mu}-m_{0}\right) \psi(x)+\mathcal{L}_{i n t} \equiv \int \mathrm{~d}^{4} x \mathcal{L} \tag{7.293}
\end{equation*}
$$

with

$$
\begin{equation*}
\mathcal{L}=\bar{\psi}(x) i \gamma^{0} \partial_{0} \psi(x)+\bar{\psi}(x)\left(i \gamma^{i} \partial_{i}-m_{0}\right) \psi(x)+\mathcal{L}_{i n t} \tag{7.294}
\end{equation*}
$$

and $\psi \equiv \psi^{A}, \bar{\psi} \equiv \psi_{A}^{\dagger} \gamma^{0},\left(\gamma^{0}\right)^{2}=\mathbb{1}$. Note that the Hamiltonian is

$$
\begin{equation*}
H=-\bar{\psi}(x)\left(i \gamma^{i} \partial_{i}-m_{0}\right) \psi(x)-\mathcal{L}_{i n t} \tag{7.295}
\end{equation*}
$$

Thus we immediately generalise the above result to field theory as

$$
\begin{equation*}
\left\langle\psi_{F}\left(\vec{x}_{F}\right), t_{F} \mid \psi_{I}\left(\vec{x}_{I}\right), t_{I}\right\rangle=\int \mathcal{D} \bar{\psi}(t, \vec{x}) \mathcal{D} \psi(t, \vec{x}) e^{i \int_{t_{I}}^{t_{F}} \mathrm{~d}^{4} x \mathcal{L}(\psi, \vec{\psi})} \tag{7.296}
\end{equation*}
$$

Note that we can integrate over $\bar{\psi}$ as opposed to $\psi^{*}$ without cost. To project initial and final states to the vacuum $|\Omega\rangle$, the same trick of replacing

$$
\begin{equation*}
m_{0} \rightarrow m_{0}-i \epsilon \tag{7.297}
\end{equation*}
$$

works with $t_{I} \rightarrow-\infty$ and $t_{F} \rightarrow \infty$. With this understanding we define the generating functional for fermionic correlation functions as

$$
\begin{align*}
Z[\eta, \bar{\eta}] & =\langle\Omega| T e^{i \int \mathrm{~d}^{4} x[\mathcal{L}(\psi, \bar{\psi})+\bar{\eta}(x) \psi(x)+\bar{\psi}(x) \eta(x)]}|\Omega\rangle \\
& =\int \mathcal{D} \bar{\psi} \mathcal{D} \psi e^{i S[\psi, \bar{\psi}]+i \bar{\eta} \cdot \psi+i \bar{\psi} \cdot \eta} \tag{7.298}
\end{align*}
$$

where the external sources $\eta, \bar{\eta}$ are Grassmann-valued classical fields. Then for instance

$$
\begin{equation*}
\langle\Omega| T \psi\left(x_{1}\right) \bar{\psi}\left(x_{2}\right)|\Omega\rangle=\left.\frac{1}{Z[0,0]} \frac{\delta}{i \delta \bar{\eta}\left(x_{1}\right)} \frac{-\delta}{i \delta \eta\left(x_{2}\right)} Z[\eta, \bar{\eta}]\right|_{\eta=0=\bar{\eta}}, \tag{7.299}
\end{equation*}
$$

where special attention must be paid to the various minus signs.
The same methods as in the bosonic theory can be applied to deduce from this the fermionic Feynman rules. For example, as will be derived in the tutorials, in the free theory completing the square yields

$$
\begin{equation*}
Z_{0}[\eta, \bar{\eta}]=Z_{0}[0,0] e^{-\bar{\eta} S_{F} \eta} \tag{7.300}
\end{equation*}
$$

where the fermionic Feynman propagator

$$
\begin{equation*}
S_{F}(x-y)=\int \frac{\mathrm{d}^{4} p}{(2 \pi)^{4}} \frac{i\left(\gamma \cdot p+m_{0}-i \epsilon\right)}{p^{2}-m_{0}^{2}+i \epsilon} e^{-i p(x-y)} \tag{7.301}
\end{equation*}
$$

is a solution to

$$
\begin{equation*}
\left(i \gamma \cdot \partial_{x}-m_{0}+i \epsilon\right)\left(-i S_{F}(x-y)\right)=\delta(x-y), \quad \text { i.e. } \quad\left(i \gamma \cdot \partial-m_{0}+i \epsilon\right)=i S_{F}^{-1} . \tag{7.302}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\langle 0| T \psi(x) \bar{\psi}(y)|0\rangle=S_{F}(x-y) \tag{7.303}
\end{equation*}
$$

Similarly one confirms the fermionic Feynman rules as stated in QFT I.

## Chapter 8

## Renormalisation of Quantum Field Theory

In QFT I we claimed without proof the following properties of Quantum Electrodynamics (QED):

- The structure of UV divergences is encoded in the following three types of diagrams:


- The four constants multiplying the divergences in these diagrams can be absorbed order by order in perturbation theory in the counterterms of the renormalized Lagrangian such that all physical amplitudes are finite. For this recall chapter 5.9 of QFT I.

While we demonstrated the second point at 1-loop order in QED, a number of open questions remain:

1. Why are only the above diagrams relevant as far as the structure of UV divergences is concerned, and how do these diagrams generalise to other QFTs?
2. How does one prove finiteness of the renormalised theory beyond 1-loop order in perturbation theory?
3. What is the physics behind the renormalisation procedure?

### 8.1 Superficial divergence and power counting

In this section we address the question which diagrams are the relevant ones concerning the structure of UV divergences in QFT. We will do so for definiteness in the context of a scalar theory in $d$
dimensions with bare Lagrangian

$$
\begin{equation*}
\mathcal{L}_{0}=\frac{1}{2}\left(\partial_{\mu} \phi\right)^{2}-\frac{1}{2} m_{0}^{2} \phi^{2}-\frac{\lambda_{0}^{(n)}}{n!} \phi^{n} . \tag{8.1}
\end{equation*}
$$

Consider a connected diagram with $E$ external lines, $L$ loops, $V$ vertices and $I$ internal propagators. We had seen that for such diagrams

$$
\begin{equation*}
L=I-V+1 . \tag{8.2}
\end{equation*}
$$

Since each external line ends with 1 end at a vertex and each internal line with 2 ends, we have

$$
\begin{equation*}
n V=2 I+E \tag{8.3}
\end{equation*}
$$

Now, naively the structure of momentum integrals in the UV for such a diagram is

$$
\begin{equation*}
\frac{\int \mathrm{d}^{d} k_{1} \ldots \int \mathrm{~d}^{d} k_{L}}{k_{1}^{2} \ldots k_{I}^{2}} \tag{8.4}
\end{equation*}
$$

since each loop gives a momentum integral $\int \mathrm{d}^{d} k$ and each internal line contributes a propagator scaling like $\left(k^{2}-m^{2}\right)^{-1}$. Let us define $D$ to be the difference of the power of the momentum in nominator and denominator. In the above theory

$$
\begin{equation*}
D=d L-2 I . \tag{8.5}
\end{equation*}
$$

If we introduce a momentum cutoff in the UV by replacing

$$
\begin{equation*}
\int^{\infty} \mathrm{d} k \rightarrow \int^{\Lambda} \mathrm{d} k \tag{8.6}
\end{equation*}
$$

then according to this naive analysis the amplitude $\mathcal{M}$ for the process associated with a diagram scales as follows with $\Lambda$ :

$$
\begin{array}{ll}
D>0: \mathcal{M} \sim \Lambda^{D} & \text { (superfically divergent) } \\
D<0: \mathcal{M} \sim \Lambda^{-|D|} & \text { (superficially finite) } \\
D=0: \mathcal{M} \sim \log (\Lambda) & \text { (superficially log-divergent). }
\end{array}
$$

According to this reasoning, as $\Lambda \rightarrow \infty$ only diagrams with $D \geq 0$ are divergent. Therefore $D$ is called superficial degree of divergence. The term superficial indicates that $D$ does not always reflect the actual divergence or finiteness properties of a diagram, due to one of the three possible exceptions:

- If $D \geq 0$, then a diagram may still be actually finite if a sufficient amount of symmetry constrains the form of the amplitude or leads to cancellations among infinite terms. For example, supersymmetry leads to such cancellations, and the theory with maximal supersymmetry, $N=4$ Super-Yang-Mills theory in 4D, is even finite to all orders.
- If $D<0$, a diagram may be actually divergent if the divergence is due to a simpler divergent subdiagram.
- Tree-level diagrams have $D=0$, but are finite.

To illustrate this consider $\phi^{4}$ theory in $d=4$.

1. The following diagram has $D=0$ and is superficially and actually logarithmically divergent, as we will see later.

2. The following diagram has $D=-2$ and is superficially and actually finite.

3. The following diagram still has $D=-2$, but is actually divergent. However, a closer look reveals that the divergence is due to diagram 1, which appears as a divergent subdiagram with $D=0$.


Conclusions: Except for the issue of divergent sudiagrams, all divergences come at most from the superficially divergent diagrams.

Given its importance, it is useful to derive an explicit expression for $D$. In the theory under consideration ( $\phi^{n}$ in $d$ dimensions) we have

$$
\begin{equation*}
D=d L-2 I=d(I-V+1)-2 I \tag{8.7}
\end{equation*}
$$

with $I=\frac{1}{2}(n V-E)$ and thus

$$
\begin{equation*}
D=d-\left(d-n \frac{d-2}{2}\right) V-\frac{d-2}{2} E \text {. } \tag{8.8}
\end{equation*}
$$

Equation (8.8) can also be obtained by dimensional analysis: The action

$$
\begin{equation*}
S=\int \mathrm{d}^{d} x\left(-\frac{1}{2} \phi\left(\partial^{2}+m_{0}^{2}\right) \phi-\frac{\lambda_{0}^{(n)}}{n!} \phi^{n}\right) . \tag{8.9}
\end{equation*}
$$

has mass dimension zero (in natural units with $\hbar=1=c$ ) and therefore

$$
\begin{equation*}
[\phi]=\frac{d-2}{2}, \quad\left[\lambda_{0}^{(n)}\right]=-n \frac{d-2}{2}+d . \tag{8.10}
\end{equation*}
$$

Consider a diagram with $E$ external lines. To determine the mass dimension of the associated amplitude, we note that it would receive a tree-level contribution $\mathcal{M}^{(E)}=-i \lambda^{(E)}$ if there were a term $\lambda^{(E)} \phi^{E}$ in $\mathcal{L}$. Thus

$$
\begin{equation*}
\left[\mathcal{M}^{(E)}\right]=\left[\lambda^{(E)}\right]=d-E \frac{d-2}{2} . \tag{8.11}
\end{equation*}
$$

Furthermore, if $\mathcal{M}^{(E)}$ has superficial degree $D$, then

$$
\begin{equation*}
\mathcal{M}^{(E)} \sim\left(\lambda_{o}^{(n)}\right)^{V} \Lambda^{D} \tag{8.12}
\end{equation*}
$$

unless one of the 3 exceptions is at work, where $V$ is the number of $n$-point vertices. Therefore

$$
\begin{equation*}
V\left[\lambda_{0}^{(n)}\right]+D=d-E \frac{d-2}{2} \tag{8.13}
\end{equation*}
$$

and thus

$$
\begin{equation*}
D=d-\left(d-n \frac{d-2}{2}\right) V-\frac{d-2}{2} E=d-\left[\lambda_{0}^{(n)}\right] V-\frac{d-2}{2} E . \tag{8.14}
\end{equation*}
$$

The appearance of the mass dimension of the coupling constant is very important as we will see in a second.

### 8.2 Renormalisability and BPHZ theorem

The UV properties of a theory depend crucially on the sign of the coefficient of $V$ in the expression (8.8) for $D$. To see this consider the following examples:

1. For $\phi^{4}$ theory in 4 dimensions $(d=4, n=4)$ we find $D=4-E$, i.e. $D$ depends only on $E$, not on $V$ or $L$. There exists only a finite number of superficially divergent amplitudes. Since the action is invariant under $\phi \rightarrow-\phi$, any diagram with an odd number of external lines vanishes, and the only superficially divergent amplitudes are the following:


Since $D$ is independet of $L$, superficial divergences appear for these amplitudes at every order in perturbation theory.
2. In $\phi^{3}$ theory in 4 dimensions $(d=4, n=3)$, we have $D=4-V-E$. Now only a finite number of superficially divergent Feynman diagrams exist since $V$ appears directly and with a negative sign.
3. In $\phi^{5}$ theory in 4 dimensions ( $d=4, n=5$ ), we have $D=4+V-E$. There exists an infinite number of superficially divergent amplitudes since for every value of $E$ diagrams with high enough $V$ diverge.

Definition: A theory is called

- (power-counting) renormalisable if the number of superficially divergent amplitudes is finite, but superficial divergences appear at every order in perturbation theory.
- (power-counting) super-renormalisable if the number of superficially divergent Feynman diagrams is finite.
- (power-counting) non-renormalisable if the number of superficially divergent amplitudes is infinite.

From $D=d-\left[\lambda^{(n)}\right] V-\frac{d-2}{2} E$ (and analogous formulae for other theories) we draw the following important conclusion:
A theory is

- renormalisable if its coupling has vanishing mass dimension.
- super-renormalisable if its coupling has positive mass dimension.
- non-renormalisable if its coupling has negative mass dimension.

The significance of the concept of (power-counting) renormalisablity is summarised in the BPHZ theorem ${ }^{1}$

- Ignore the issue of divergent subgraphs for a moment. Then if a theory is (power-counting) renormalisable, at each order in perturbation theory only a finite number of divergent diagrams, parametrized by a finite number of divergent constants, appear. One can absorb these divergences order by order in the counterterms of the renormalized Lagrangian such that all physical amplitudes are finite.
- The counterterms create new Feynman diagrams relevant at the next order. These will cancel the divergences of the divergent subdiagrams (if present) at the next order.

[^39]- All of this leads to a perturbative adjustment of a finite number of counterterms in the renormalised Lagrangian, and thus predictivity is maintained.

The non-trivial aspect of this theorem concerns the complete cancellation of divergent subdiagrams by lower-order counterterms. Rather than presenting a full proof, which can be found in Zimmermann's 1970 article (see also Kugo's book), we will demonstrate this in the next section in the context of renormalisation at two-loop order. For simplicity we do this for $\phi^{4}$ theory.

### 8.3 Renormalisation of $\phi^{4}$ theory up to 2-loops

We follow mostly the procedure of Peskin-Schröder, Chapters 10.2 and 10.5. Our starting point is the bare Lagrangian

$$
\begin{equation*}
\mathcal{L}=\frac{1}{2}\left(\partial_{\mu} \phi\right)^{2}-\frac{1}{2} m_{0}^{2} \phi^{2}-\frac{\lambda_{0}}{4!} \phi^{4} \tag{8.15}
\end{equation*}
$$

in 4 dimensions with $\left[\lambda_{0}\right]=0$. As noted already above the superficially divergent diagrams are given by

- the vacuum energy with $E=0, D=4$, which is trivially absorbable into $V_{0}$ :
- the propagator with $E=2, D=2$ :

- the 4-point amplitude with $E=4, D=0$ :


Now comes a simple, but important observation regarding dimensional analysis that we will make frequent use of: In a theory with a dimensionless coupling constant the dimension of the amplitude agrees with $D$, see eq. (8.12). We can therefore parametrise the second amplitude with $D=2$ in terms of the momentum cutoff $\Lambda$ as

$$
\begin{equation*}
p \longrightarrow p=a \Lambda^{2}+b p \Lambda+c p^{2} \log (\Lambda)+\text { finite terms. } \tag{8.16}
\end{equation*}
$$

This is to be viewed as a Taylor expansion in $p$ with dimensionless constants $a, b, c$. Note that due to the symmetry $p \rightarrow-p$ of $\mathcal{L}$ in momentum space $b \equiv 0$.
By the same logic, the 4-point amplitude with $D=0$ can be parametrized (with $d$ dimensionless) as


Definition: If the divergence is multiplied by a simple monomial in $p$ as above, it is called local divergence. In the presence of divergent subgraphs more complicated expressions of $p$ will appear. The BPHZ theorem states that it is possible to absorb the 3 divergent constants $a, c, d$ into 3 counterterms in the renormalised Lagrangian, order by order in renormalized perturbation theory. As discussed in the context of QED in QFT1 we define the renormalised field

$$
\begin{equation*}
\phi_{r}=Z^{-\frac{1}{2}} \phi \tag{8.17}
\end{equation*}
$$

with $Z$ the wavefunction renormalisation and rewrite $\mathcal{L}$ as

$$
\begin{equation*}
\mathcal{L}=\frac{1}{2}\left(\partial_{\mu} \phi_{r}\right)^{2}-\frac{1}{2} m^{2} \phi_{r}^{2}-\frac{\lambda}{4!} \phi_{r}^{4}+\frac{1}{2} \delta_{Z}\left(\partial_{\mu} \phi_{r}\right)^{2}-\frac{1}{2} \delta_{m} \phi_{r}^{2}-\frac{\delta_{\lambda}}{4!} \phi_{r}^{4} \tag{8.18}
\end{equation*}
$$

where $\delta_{Z}=Z-1, \delta_{m}=m_{0}^{2} Z-m^{2}$ and $\delta_{\lambda}=\lambda_{0} Z^{2}-\lambda$. We henceforth compute with the renormalised Feynman rules

and no factors of $\sqrt{Z}$ for in- and out-going states in the computation of scattering amplitudes. To define the meaning of $m$ and $\lambda$ we impose the following renormalisation conditions:


Here $s, t, u$ are the Mandelstam variables, which read with the above convention for the arrow of the momenta,

$$
\begin{equation*}
s=\left(p_{1}+p_{2}\right)^{2}, \quad t=\left(p_{1}-p_{3}\right)^{2}, \quad u=\left(p_{1}-p_{4}\right)^{2} . \tag{8.19}
\end{equation*}
$$

Now

$$
\begin{equation*}
-=\frac{i}{p^{2}-m^{2}-M^{2}\left(p^{2}\right)}, \tag{8.20}
\end{equation*}
$$

where $-i M^{2}\left(p^{2}\right)$ is the amplitude of the $1 P I$ amputated diagrams. Thus the first renormalisation condition implies:

$$
\begin{equation*}
\left.M^{2}\left(p^{2}\right)\right|_{p^{2}=m^{2}} \stackrel{!}{=} 0,\left.\quad \frac{\mathrm{~d}}{\mathrm{~d} p^{2}} M^{2}\left(p^{2}\right)\right|_{p^{2}=m^{2}} \stackrel{!}{=} 0 . \tag{8.21}
\end{equation*}
$$

Note that the renormalisation conditions are arbitrary and define the renormalisation scheme.

### 8.3.1 1-loop renormalisation

At 1-loop order

and thus ${ }^{2}$

$$
\begin{equation*}
\ldots(-i \lambda)-\frac{1}{2} \int \frac{\mathrm{~d}^{4} k}{(2 \pi)^{4}} \frac{i}{k^{2}-m^{2}+i \epsilon}+i\left(p^{2} \delta_{Z}^{(1)}-\delta_{m}^{(1)}\right) \equiv-\left.i M^{2}\left(p^{2}\right)\right|_{1-\mathrm{loop}} . \tag{8.22}
\end{equation*}
$$

We compute the integral by dimensional regularisation. The result (see exercise sheet 5 for a derivation) is

$$
\begin{equation*}
\int \frac{\mathrm{d}^{d} k}{(2 \pi)^{d}} \frac{i}{k^{2}-m^{2}}=\frac{1}{(4 \pi)^{d / 2}} \frac{\Gamma\left(1-\frac{d}{2}\right)}{\left(m^{2}\right)^{1-d / 2}} . \tag{8.23}
\end{equation*}
$$

For $d=4-\epsilon$ and with the help of

$$
\begin{equation*}
\Gamma\left(-1+\frac{\epsilon}{2}\right)=-\frac{2}{\epsilon} e^{-\frac{1}{2} \gamma \cdot \epsilon}\left(1+\frac{\epsilon}{2}\right)+\boldsymbol{O}\left(\epsilon^{3}\right) \tag{8.24}
\end{equation*}
$$

we get

$$
\bigcirc=i \frac{\lambda}{2}\left(1+\frac{2}{\epsilon}\right) \frac{\left(m^{2}\right)^{1-\epsilon / 2}}{(4 \pi)^{2-\epsilon / 2}} e^{-\frac{1}{2} \gamma \epsilon}=\frac{i}{2} \frac{\lambda}{16 \pi^{2}} m^{2}\left(\frac{2}{\epsilon}+1-\log \left(m^{2}\right)\right)\left(e^{-\gamma} 4 \pi\right)^{\frac{1}{2} \epsilon},
$$

where we used

$$
\begin{equation*}
\left(m^{2}\right)^{-\epsilon / 2}=e^{-\frac{\epsilon}{2} m^{2}}=1-\frac{\epsilon}{2} \log \left(m^{2}\right)+O\left(\epsilon^{3}\right) . \tag{8.25}
\end{equation*}
$$

Note that the expected quadratic divergence $a \Lambda^{2}$ corresponds, in dimensional regularization, to the divergent term

$$
\begin{equation*}
\frac{i \lambda}{16 \pi^{2}} m^{2} \frac{1}{\epsilon} \tag{8.26}
\end{equation*}
$$

More precisely, for the integral with a UV momentum cutoff, the divergence would be proportional to

$$
\begin{equation*}
\Lambda^{2}-m^{2} \log \left(\frac{m^{2}+\Lambda^{2}}{m^{2}}\right) \tag{8.27}
\end{equation*}
$$

As a rule of thumb, in a theory with dimensionless coupling constant, the order in $\Lambda$ corresponding to a $\frac{1}{\epsilon}$-pole in dimensional regularization agrees with the mass dimension of the prefactor: Here $\left[\lambda m^{2}\right]=2$ and thus the divergence is expected to scale as $\Lambda^{2}$.

[^40]Altogether

$$
\begin{equation*}
-i M^{2}\left(p^{2}\right)=-\frac{i \lambda}{2} \frac{1}{4 \pi^{d / 2}} \frac{\Gamma\left(1-\frac{d}{2}\right)}{\left(m^{2}\right)^{1-d / 2}}+i\left(p^{2} \delta_{Z}^{(1)}-\delta_{m}^{(1)}\right) \tag{8.28}
\end{equation*}
$$

and the renormalisation conditions imply

$$
\begin{equation*}
\left.\frac{\mathrm{d}}{\mathrm{~d} p^{2}} M^{2}\left(p^{2}\right)\right|_{p^{2}=m^{2}} \stackrel{!}{=} 0 \Rightarrow \delta_{Z}^{(1)} \stackrel{!}{=} 0 \tag{8.29}
\end{equation*}
$$

and

$$
\begin{equation*}
\left.M^{2}\left(p^{2}\right)\right|_{p^{2}=m^{2}} \stackrel{!}{=} 0 \Rightarrow \delta_{m}^{(1)} \stackrel{!}{=}-\left.\frac{\lambda}{2(4 \pi)^{d / 2}} \frac{\Gamma\left(1-\frac{d}{2}\right)}{\left(m^{2}\right)^{1-d / 2}}\right|_{d \rightarrow 4} \tag{8.30}
\end{equation*}
$$

Note that at 1-loop $M^{2}\left(p^{2}\right) \equiv 0 \forall p^{2}$ and the first non-trivial corrections to the propagator appear at the 2-loop level. This, however, is a specialty of $\phi^{4}$ theory in four dimensions.

The computation of the 4 -point amplitude at 1-loop level involves the following diagrams,


The first loop amplitude corresponds to the s-channel,

where $p=p_{1}+p_{2}$. Then

$$
\begin{equation*}
i \mathcal{M}\left(p_{1} p_{2} \rightarrow p_{3} p_{4}\right)=-i \lambda+(-i \lambda)^{2}[i V(s)+i V(t)+i V(u)]-i \delta_{\lambda}^{(1)}+O\left(\lambda^{3}\right) \tag{8.31}
\end{equation*}
$$

The renormalisation condition implies

$$
\begin{equation*}
\left.i \mathcal{M}\right|_{s=4 m^{2}, t=u=0}=-i \lambda \Rightarrow(-i \lambda)^{2}\left[i V\left(4 m^{2}\right)+2 i V(0)\right]-i \delta_{\lambda}^{(1)} \stackrel{!}{=} 0 . \tag{8.32}
\end{equation*}
$$

Therefore

$$
\begin{equation*}
\delta_{\lambda}^{(1)} \stackrel{!}{=}-\lambda^{2}\left(V\left(4 m^{2}\right)+2 V(0)\right) \equiv \delta_{\lambda}^{(1)}(s)+\delta_{\lambda}^{(1)}(u+t) . \tag{8.33}
\end{equation*}
$$

The computation of $V\left(p^{2}\right)$ proceeds again via dimensional regularisation and is performed on exercise sheet 5 , with the result

$$
\begin{align*}
V\left(p^{2}\right) & =-\frac{1}{2} \frac{\Gamma(\epsilon}{(4 \pi)^{2-\epsilon / 2}} \int_{0}^{1} \mathrm{~d} x \frac{1}{\left[m^{2}-x(1-x) p^{2}\right] \epsilon}  \tag{8.34}\\
& =-\frac{1}{32 \pi^{2}} \int_{0}^{1} \mathrm{~d} x\left[\frac{2}{\epsilon}-\gamma+\log (4 \pi)-\log \left(m^{2}-x(1-x) p^{2}\right]\right.
\end{align*}
$$

Then

$$
\begin{equation*}
\delta_{\lambda}^{(1)}=\frac{\lambda^{2}}{32 \pi^{2}} \int_{0}^{1} \mathrm{~d} x\left[\frac{6}{\epsilon}-3 \gamma+3 \log (4 \pi)-\log \left(m^{2}(1-4 x(1-x))\right)-2 \log \left(m^{2}\right)\right] . \tag{8.35}
\end{equation*}
$$

Computing $i \mathcal{M}\left(p_{1} p_{2} \rightarrow p_{3} p_{4}\right)$ up to $O\left(\lambda^{2}\right)$ gives:

$$
\begin{align*}
i \mathcal{M}\left(p_{1} p_{2} \rightarrow p_{3} p_{4}\right)= & -i \lambda-\frac{i \lambda^{2}}{32 \pi^{2}} \int_{0}^{1} \mathrm{~d} x\left[\log \left(\frac{m^{2}-x(1-x) s}{m^{2}-x(1-x) 4 m^{2}}\right)\right.  \tag{8.36}\\
& \left.+\log \left(\frac{m^{2}-x(1-x) t}{m^{2}}\right)+\log \left(\frac{m^{2}-x(1-x) u}{m^{2}}\right)\right] .
\end{align*}
$$

This is manifestly finite and depends on two measurable input parameters $m$ and $\lambda$.

### 8.3.2 Renormalisation at 2-loop

We now move on to 2-loop order and demonstrate that all divergences can be absorbed in 2-loop counterterms. In particular we need to show that divergences due to divergent subdiagrams are cancelled by the 1-loop counterterms that we have already included.

To exemplify this aspect of the BPHZ theorem we consider

at 2-loop-level.

1. In diagrams of the type

the 1-loop counterterm to the propagator in the second diagram cancels the extra divergence in the first diagram. In fact, as a speciality of $\phi^{4}$-theory, we have seen that, at 1-loop level,


So the above diagrams give no contribution at order $\lambda^{3}$.
2. There are $5 \times 3+1$ diagrams of the form

plus the corresponding t-channel diagrams

plus the u-channel diagrams

as well as the 2-loop counterterm, which we split as


We consider only the s-channel, since $t$-and u-channel work completely analogously. The important point to notice is that all momentum independent divergences can be absorbed into


One must therefore show that the momentum dependent divergences cancel among the five schannel diagrams. Following Peskin-Schröder, they can be grouped into
group 1:

group 2:

and group 3:


Note that group 1 contains only the s-channel 1-loop counter terms, while in group 2 and 3 the combined $(t+u)$-channel counter terms appear. The divergences of the diagrams of group 1 are very easy to deal with: In terms of our previous definition

with $p=p_{1}+p_{2}=p_{3}+p_{4}$ we have

and


For the counter term diagram we used our result that $-i \delta_{\lambda}^{(1)}(s)=\lambda^{2}\left(i V\left(4 m^{2}\right)\right)$. Adding both contributions up and completing the square leads to

$$
\begin{equation*}
(-i \lambda)^{3}\left(-\left(V\left(p^{2}\right)-V\left(4 m^{2}\right)\right)^{2}+\left(V\left(4 m^{2}\right)\right)^{2}\right) \tag{8.37}
\end{equation*}
$$

for group 1. From the explicit expression of $V\left(p^{2}\right)$ we deduce that

$$
\begin{equation*}
V\left(p^{2}\right)-V\left(4 m^{2}\right)=\text { finite }, \tag{8.38}
\end{equation*}
$$

while $\left(V\left(4 m^{2}\right)\right)^{2}$ has merely a momentum-independent divergence, which can be absorbed by the 2-loop counter term


Note that $\sim\left(\frac{2}{\epsilon}\right)^{2}$ contains a divergence of the form $\frac{2}{\epsilon}$. More generally, higher loop-order divergences give higher powers in $1 / \epsilon$.

Now consider the diagram from group 2, given by


$$
\begin{equation*}
=(-i \lambda)^{3} \int \frac{\mathrm{~d}^{4} k}{(2 \pi)^{4}} \frac{i}{k^{2}-m^{2}} \frac{i}{(k+p)^{2}-m^{2}} i V\left(\left(p_{3}+k\right)^{2}\right) \tag{8.39}
\end{equation*}
$$

with $p=p_{1}+p_{2}=p_{3}+p_{4}$. This diagram exemplifies the appearance of a nested singularity, which appears when 2 divergent subdiagrams share a common propagator. In this case, so-called non-local divergences appear, for which the prefactor of the singularity is not a simple polynomial in $p$, but a more complicated function in $p$.

Integrals of this kind would keep us from the street for a while. We do not explicitly perform the computation, which is straightforward, but a little lengthy, and merely quote the result from Peskin-Schröder, Chapter 10: When the dust has settled, the amplitude can be written as the sum of 2 terms:

- Term 1 has only a momentum independent singularity, which can be absorbed into the 2-loop counter-term.
- Term 2 is given by

$$
\begin{equation*}
-\frac{i \lambda^{3}}{2(4 \pi)^{4}} \frac{2}{\epsilon} \int_{0}^{1} \mathrm{~d} y\left[\frac{1}{\epsilon}-\gamma+\log (4 \pi)-\log \left(m^{2}-y(1-y) p^{2}\right)\right] \tag{8.40}
\end{equation*}
$$

and exhibits 2 types of singularities:
The double pole $\frac{2}{\epsilon} \frac{1}{\epsilon}$ indicates the higher-loop nature of the singularity, as before. It is due to the region in momentum space where both loops diverge. By contrast, the term

$$
\begin{equation*}
\frac{2}{\epsilon} \log \left(m^{2}-y(1-y) p^{2}\right) \tag{8.41}
\end{equation*}
$$

is the non-local singularity advertised above. The factor of $2 / \epsilon$ is due to the regime in momentum space where only 1 of the subdiagrams diverges, and the $\log \left(p^{2}\right)$ represents the finite QFT process around this from the other subdiagram.
Explicit computation confirms that these singularities are cancelled by the counter-term diagram


This completes the proof that at 2-loop order, all divergences, including those of divergent subdiagrams, can be absorbed into the counter-terms. A formalization of this procedure leads to the BPHZ theorem. More on the ideas behind its proof can be found e.g. in Kugo.

### 8.4 Renormalisation of QED revisited

In QED with interaction $\sim \bar{\psi} \gamma^{\mu} A_{\mu} \psi$, the superfical degree of divergence of an amplitude with

$$
\begin{align*}
& E_{\gamma}=\text { number of external photons }  \tag{8.42}\\
& E_{e}=\text { number of external electrons }
\end{align*}
$$

is (see Exercise Sheet 5 for a derivation)

$$
\begin{equation*}
D=4-E_{\gamma}-\frac{3}{2} E_{e} \tag{8.43}
\end{equation*}
$$

In particular, since $[e]=0$ the theory is renormalisable. The following amplitudes have $D \geq 0$ :

1. The vacuum energy with $D=4$;
2. The photon propagator $\sim \sim$ with $D=2$, the electron propagator $\longrightarrow$ with $D=1$ and the $D=0$ vertex

3. The 1-photon amplitude $\sim$ with $D=3$;
4. The 3-photon amplitude with $D=1$

5. The 4-photon amplitude with $D=0$


Due to both discrete and continuous symmetries of the theory, only the diagrams from groups 1. and 2. are actually divergent.

- Amplitudes 3. and 4. vanish to all orders due to the discrete symmetry of charge conjugation:
$\mathcal{L}_{\mathrm{QED}}$ is invariant under

$$
\begin{equation*}
j^{\mu} \rightarrow-j^{\mu} \quad \text { and } A^{\mu} \rightarrow-A^{\mu} \tag{8.44}
\end{equation*}
$$

with $j^{\mu} \sim \bar{\psi} \gamma^{\mu} \psi$. Now, each external photon couples via its current $j^{\mu}$, e.g.

$$
\sim \sim \sim \sim \sim\langle\Omega| T j^{\mu}|\Omega\rangle
$$

Thus a diagram with only an odd number of external photons vanishes.

- Diagram 5. is non-zero, but actually finite as a consequence of gauge symmetry. This can be shown by exploiting the Ward identities.
Note that diagram 5. is responsible for the non-linearity of QED due to the scattering of photons with each other induced by loop effects.

Thus as claimed in QFT I, as far as the UV divergences are concerned it suffices to consider the diagrams 2. (the contribution of 1 . is easily absorbed into the vacuum energy term). But again symmetry lowers the actual degree of divergence for the first two diagrams of 2. to logarithmic ones:

- By dimensional analysis, the fermion propagator, having $D=1$, can be Taylor-expanded as

$$
\begin{equation*}
\xrightarrow{p}=A_{0}+A_{1} p+A_{2} p^{2}+\ldots \tag{8.45}
\end{equation*}
$$

with

$$
\begin{equation*}
A_{0}=a_{0} \Lambda, \quad A_{1}=a_{1} \log \Lambda, \quad A_{2} \text { and higher: finite as } \Lambda \rightarrow \infty \tag{8.46}
\end{equation*}
$$

and $a_{0}, a_{1}$ dimensionless. But then the mass shift of the electron would be proportional to $\Lambda$. Instead this shift must be proportional to the electron mass $m_{e}$ itself. This is because if $m_{e} \equiv 0$, then the theory possesses chiral or axial symmetry, because the Lagrangian can be written as

$$
\begin{equation*}
\mathcal{L}=\bar{\psi}_{L} i \gamma \cdot \partial \psi_{L}+\bar{\psi}_{R} i \gamma \cdot \partial \psi_{R}-e j^{\mu} A_{\mu} \tag{8.47}
\end{equation*}
$$

and is invariant under independent rotations of the left- and right-chiral spinors

$$
\begin{equation*}
\psi_{L} \rightarrow e^{i \alpha} \psi_{L}, \quad \psi_{R} \rightarrow e^{i \beta} \psi_{R} \tag{8.48}
\end{equation*}
$$

This symmetry, however, forbids a mass term, which must be of the form $m_{e} \bar{\psi}_{R} \psi_{L}$. Thus $A_{0}$ must contain an explicit power of $m_{e}$. By dimensional analysis there is no more room for a factor of $\Lambda$, and $A_{0}$ can at worst be of the form

$$
\begin{equation*}
A_{0}=\tilde{a}_{0} m_{e} \log \Lambda, \tag{8.49}
\end{equation*}
$$

as is in fact confirmed by computation.

- The photon propagator is, by dimensional analysis, of the form
with

$$
\begin{equation*}
B_{0}=b_{0} \Lambda^{2}, \quad B_{1}^{\mu}=b_{1}^{\mu} \Lambda, \quad B_{2}=b_{2} \log \Lambda \tag{8.51}
\end{equation*}
$$

and $b_{i}$ dimensionless. But by the Ward identities we know that

$$
\begin{equation*}
\mu \sim \sim \sim v=\left(\eta^{\mu v} q^{2}-q^{\mu} q^{v}\right) \Pi\left(q^{2}\right) \tag{8.52}
\end{equation*}
$$

and thus $b_{0}=b_{1}=0$.

Thus all diagrams of 2 . are logarithmically divergent. Renormalisation proceeds as discussed, at 1-loop level, in QFT I.

### 8.5 The renormalisation scale

Renormalisation automatically introduces a mass scale $\mu$ - the renormalisation scale - in the quantum theory via the renormalisation conditions. In our renormalisation of $\phi^{4}$ theory we have so far chosen to identify $\mu$ with the mass $m$ of $\phi$ because our renormalisation conditions have been

$$
\begin{equation*}
\left.M^{2}\left(p^{2}\right)\right|_{p^{2}=m^{2}}=0=\left.\frac{\mathrm{d}}{\mathrm{~d} p^{2}} M^{2}\left(p^{2}\right)\right|_{p^{2}=m^{2}} \tag{8.53}
\end{equation*}
$$

and


More generally, the renormalisation scale $\mu$ can be an independent scale since the renormalisation conditions are arbitrary. Since this observation is so crucial for all that follows we would like to illustrate it very explicitly. To this end let us consider the massless theory

$$
\begin{equation*}
\mathcal{L}=\frac{1}{2}(\partial \phi)^{2}-\frac{\lambda_{0}}{4!} \phi_{0}^{4} \tag{8.54}
\end{equation*}
$$

in $\mathrm{d}=4$. Note that this theory has no classical scale since $m_{0}=0$. However, as we will now see, the scale-invariance of the classical theory is broken in the quantum theory due to the appearance of the renormalisation scale. For instance, we can impose the renormalisation conditions

1. $M^{2}\left(p^{2}\right)=0$ at $p^{2}=-\mu^{2}$, i.e. for spacelike momenta,
2. $\frac{\mathrm{d}}{\mathrm{d} p^{2}} M^{2}\left(p^{2}\right)=0$ at $p^{2}=-\mu^{2}$,
3. and for the vertex:


With $\phi=Z^{-1 / 2} \phi_{0}$, conditions 1. and 2. imply

$$
\begin{equation*}
\left.\int \mathrm{d}^{4} x e^{i p \cdot x}\langle\Omega| T \phi(x) \phi(0)|\Omega\rangle\right|_{p^{2}=-\mu^{2}}=\left.\frac{i}{p^{2}}\right|_{p^{2}=-\mu^{2}}, \tag{8.55}
\end{equation*}
$$

i.e.

$$
\begin{equation*}
\left.\int \mathrm{d}^{4} x e^{i p \cdot x}\langle\Omega| T \phi_{0}(x) \phi_{0}(0)|\Omega\rangle\right|_{p^{2}=-\mu^{2}}=\left.\frac{i Z}{p^{2}}\right|_{p^{2}=-\mu^{2}} \tag{8.56}
\end{equation*}
$$

In particular, the wavefunction renormalisation factor $Z$ used in $\phi=Z^{-1 / 2} \phi_{0}$ is not the residue of the propagator at the physical mass $p^{2}=0$, but has the meaning as given in the previous equation.
Now, to fully appreciate the appearance of the renormalisation scale $\mu$, consider the renormalised connected 4-correlator

$$
\begin{equation*}
G^{(4)}\left(x_{1}, . ., x_{4}\right)=\left.\langle\Omega| T \phi\left(x_{1}\right) \ldots \phi\left(x_{4}\right)|\Omega\rangle\right|_{\text {connected }} . \tag{8.57}
\end{equation*}
$$

Its Fourier transformation is

$$
\begin{equation*}
G^{(4)}\left(p_{1}, . ., p_{4}\right)=\left[-i \lambda+(-i \lambda)^{2}(i V(s)+i V(t)+i V(u))-i \delta_{\lambda}+O\left(\lambda^{3}\right)\right] \prod_{i=1}^{4} \frac{i}{p_{i}^{2}} \tag{8.58}
\end{equation*}
$$

Note the product $\prod_{i=1}^{4} \frac{i}{p_{i}^{2}}$ from the external legs because we are computing the 4-point correlator, not the scattering amplitude $i \mathcal{M}\left(p_{1} p_{2} \rightarrow p_{3} p_{4}\right)$. Renormalisation condition 3. implies

$$
\begin{align*}
\delta_{\lambda} & =(-i \lambda)^{2} \cdot 3 V\left(-\mu^{2}\right)=\frac{3 \lambda^{2}}{2(4 \pi)^{d / 2}} \int_{0}^{1} \mathrm{~d} x \frac{\Gamma\left(2-\frac{d}{2}\right)}{\left(x(1-x) \mu^{2}\right)^{2-d / 2}}  \tag{8.59}\\
& =\frac{3 \lambda^{2}}{2(4 \pi)^{2}}\left[\frac{2}{\epsilon}-\log \mu^{2}+\text { finite terms }\right]
\end{align*}
$$

in dimensional regularization, whereas in cutoff regularisation the term in brackets would be of the form

$$
\begin{equation*}
\log \frac{\Lambda^{2}}{\mu^{2}}+\text { finite. } \tag{8.60}
\end{equation*}
$$

The crucial observation is that $G^{(4)}\left(p_{1}, \ldots, p_{4}\right)$ is finite, but

- explicitly depends on the renormalisation scale $\mu$ through the counterterm $\delta_{\lambda}$ and
- implicitly depends on $\mu$, at 1-loop, via $\lambda$ because $\lambda$ is defined via renormalisation condition 3 . at the scale $\mu$.
(At higher loop, there will also be an implicit $\mu$-dependence via $Z$ as we will see below.)


### 8.6 The Callan-Symanzyk (CS) equation

The CS equation quantifies the dependence of various quantities on the renormalisation scale $\mu$. To derive it we be begin with the bare correlator (for definiteness in massive $\phi^{4}$ theory)

$$
\begin{equation*}
G_{n}^{(0)}\left(x_{1}, \ldots, x_{n}\right)=\left.\langle\Omega| T \prod_{i} \phi_{0}\left(x_{i}\right)|\Omega\rangle\right|_{\text {conn. }} \tag{8.61}
\end{equation*}
$$

computed from the bare Feynman rules in terms of the parameters $\lambda_{0}, m_{0}$ appearing in the bare Lagrangian $\mathcal{L}_{0}$ as well as the cutoff $\Lambda$ (or $\epsilon$ ). There is no explicit dependence of $G_{n}^{(0)}\left(x_{1}, \ldots, x_{n}\right)$ on the renormalisation scale $\mu$, i.e.

$$
\begin{equation*}
G_{n}^{(0)}\left(x_{1}, \ldots, x_{n}\right)=G_{n}^{(0)}\left(x ; \lambda_{0}, m_{0} ; \Lambda\right) \tag{8.62}
\end{equation*}
$$

and this quantity satisfies the equation

$$
\begin{equation*}
\left.\frac{\mathrm{d}}{\mathrm{~d} \mu} G_{n}^{(0)}\left(x ; \lambda_{0}, m_{0}\right)\right|_{\lambda_{0}, m_{0} \text { fixed }}=0 \tag{8.63}
\end{equation*}
$$

Here it is important that we do not vary with respect to $\lambda_{0}$ and $m_{0}$, but only state that no explicit dependence on $\mu$ is given. The reason we stress this is that in bare perturbation theory, $\lambda_{0}$ and $m_{0}$ must themselves depend on the cutoff $\Lambda$ in such a way as to cancel the divergences, and this definition of the bare parameters in turn involves a renormalisation scale (cf the discussion in section 5.9.1). However, this dependence on $\mu$ is only implicit via $\lambda_{0}$ and $m_{0}$, and as long as we do not include variation with respect to these bare parameters, equation (8.63) is correct.
Let us compare this with the renormalised correlator

$$
\begin{equation*}
G_{n}\left(x_{1}, \ldots, x_{n}\right)=\left.\langle\Omega| T \prod_{i} \phi\left(x_{i}\right)|\Omega\rangle\right|_{\text {conn. }} \tag{8.64}
\end{equation*}
$$

with

$$
\begin{equation*}
\phi\left(x_{i}\right)=Z^{-1 / 2} \phi_{0}\left(x_{i}\right) \tag{8.65}
\end{equation*}
$$

computed via the renormalised Feynman rules including counter-terms. As demonstrated above, $G_{n}\left(x_{i}\right)$ depends on the renormalisation parameters $\lambda, m$ and the renormalisation scale $\mu$, but not on the cutoff,

$$
\begin{equation*}
G_{n}\left(x_{1}, \ldots, x_{n}\right)=G(x ; \lambda, m ; \mu) . \tag{8.66}
\end{equation*}
$$

Now, going from the bare Lagrangian to the renormalised Lagrangian is only a reformulation of the same theory. Thus the bare and the renormalized correlators must agree up to the rescaling of fields $\phi=Z^{-1 / 2} \phi_{0}$, i.e.

$$
\begin{equation*}
G_{n}\left(x_{1}, \ldots, x_{n}\right)=Z^{-n / 2} G_{n}^{(0)}\left(x_{1},, x_{n}\right) . \tag{8.67}
\end{equation*}
$$

To avoid confusion, we stress that if we compare scattering amplitudes $i \mathcal{M}$, the amplitudes computed in bare and renormalised perturbation theory agree completely since we compensate for the rescaling of $\phi$ by omitting factors of $Z^{1 / 2}$ for external states in LSZ. This is not so in the transition from $G_{n}$ to $G_{n}^{(0)}$, and thus the extra factor $Z^{-n / 2}$ appears.

Together with (8.63) we conclude that

$$
\begin{equation*}
\left.\mu \frac{\mathrm{d}}{\mathrm{~d} \mu}\left(Z^{n / 2} G_{n}(x ; \lambda, m ; \mu)\right)\right|_{\lambda_{0}, m_{0} \text { fixed }}=0 \tag{8.68}
\end{equation*}
$$

where the extra factor of $\mu$ has been included for later convenience. This can be further evaluated with the help of the chain rule. In terms of the quantities

$$
\begin{align*}
\beta_{\lambda} & :=\left.\mu \frac{\mathrm{d} \lambda}{\mathrm{~d} \mu}\right|_{\lambda_{0}, m_{0}} \\
\beta_{m^{2}} & :=\left.\mu \frac{\mathrm{d}}{\mathrm{~d} \mu} m^{2}\right|_{\lambda_{0}, m_{0}}  \tag{8.69}\\
\gamma_{\phi} & :=\left.\frac{1}{2} \frac{\mu}{Z} \frac{\mathrm{~d} Z}{\mathrm{~d} \mu}\right|_{\lambda_{0}, m_{0}}=\left.\frac{\mu}{2} \frac{\mathrm{~d}}{\mathrm{~d} \mu} \log Z\right|_{\lambda_{0}, m_{0}}
\end{align*}
$$

we arrive at the celebrated Callan-Symanzik or renormalisation group equation

$$
\begin{equation*}
\left(\mu \frac{\partial}{\partial \mu}+\beta_{\lambda} \frac{\partial}{\partial \lambda}+\beta_{m^{2}} \frac{\partial}{\partial m^{2}}+n \cdot \gamma_{\phi}\right) G_{n}(x ; \lambda, m ; \mu)=0 \tag{8.70}
\end{equation*}
$$

The physical interpretation of the functions (8.69) will be discussed in length in a subsequent section. Suffice it here to anticipate that, for instance,

$$
\begin{equation*}
\beta_{\lambda}=\mu \frac{\mathrm{d}}{\mathrm{~d} \mu} \lambda \tag{8.71}
\end{equation*}
$$

describes how the physical coupling $\lambda$ changes as we change the energy scale at which we perform an experiment ${ }^{3}$.
The significance of the CS equation is that it allows us to compute $\beta_{\lambda}, \beta_{m^{2}}, \gamma_{\phi}$ explicitly in perturbation theory by first computing $G_{n}(x ; \lambda, m ; \mu)$ and plugging it into the CS equation. This allows for a perturbative solution for $\beta_{\lambda}, \beta_{m^{2}}, \gamma_{\phi}$, order by order in the coupling constant. Here it is important to carefully compare only terms at the same perturbative order. Even before doing this, let us state two general results:

- Since $G_{n}(x ; \lambda, m ; \mu)$ is independent of the cutoff, so must be $\beta_{\lambda}, \beta_{m^{2}}, \gamma_{\phi}$.
- While the higher order terms in $\beta_{\lambda}, \beta_{m^{2}}, \gamma_{\phi}$ will depend on the renormalisation scheme, to leading order in the coupling constant the answer turns out to be universal.

Before embarking on a more systematic analysis, consider massless $\phi^{4}$ theory in the renormalisation scheme of section (8.5):

- The 2-point function satisfies

$$
\begin{equation*}
\left[\mu \frac{\partial}{\partial \mu}+\beta_{\lambda} \frac{\partial}{\partial \lambda}+2 \gamma_{\phi}\right] G_{2}(p)=0 . \tag{8.72}
\end{equation*}
$$

At 1-loop order, which for the 2-point function is $O(\lambda)$, we have seen that as a speciality of $\phi^{4}$-theory

$$
\begin{equation*}
M^{2}\left(p^{2}\right)=0 \tag{8.73}
\end{equation*}
$$

i.e.

$$
\begin{equation*}
\left.\frac{\partial}{\partial \mu} G_{2}\right|_{1-\text { loop }}=0=\left.\frac{\partial}{\partial \lambda} G_{2}\right|_{1-\text { loop }} \tag{8.74}
\end{equation*}
$$

Thus $\gamma_{\phi}=O\left(\lambda^{2}\right)$.

[^41]- The 4 -point function is given by

$$
\begin{equation*}
G_{4}\left(p_{1},, p_{4}\right)=\left[-i \lambda+(-i \lambda)^{2}(i V(s)+i V(t)+i V(u))-i \delta_{\lambda}+O\left(\lambda^{3}\right)\right] \prod_{i} \frac{i}{p_{i}^{2}} \tag{8.75}
\end{equation*}
$$

with

$$
\begin{equation*}
\delta_{\lambda}=\frac{3 \lambda^{2}}{2(4 \pi)^{2}}\left[\frac{2}{\epsilon}-\log \mu^{2}+\text { finite }\right] . \tag{8.76}
\end{equation*}
$$

Let us evaluate the CS equation

$$
\begin{equation*}
\left[\mu \frac{\partial}{\partial \mu}+\beta_{\lambda} \frac{\partial}{\partial \lambda}+4 \gamma_{\phi}\right] G_{4}\left(p_{i}\right)=0 \tag{8.77}
\end{equation*}
$$

to order $\lambda^{2}$, which corresponds to 1 -loop for the 4-point function. First, to this order

$$
\begin{equation*}
\mu \frac{\partial}{\partial \mu} G^{(4)}=\frac{3 i \lambda^{2}}{(4 \pi)^{2}} \prod_{i} \frac{i}{p_{i}^{2}} . \tag{8.78}
\end{equation*}
$$

Since $\gamma_{\phi}=O\left(\lambda^{2}\right)$, we know that $\gamma_{\phi} G_{4}=O\left(\lambda^{3}\right)$. Finally

$$
\begin{equation*}
\beta_{\lambda} \frac{\partial}{\partial \lambda} G^{(4)}=(-i+O(\lambda)) \prod_{i} \frac{i}{p_{i}^{2}} \tag{8.79}
\end{equation*}
$$

Plugging everything into the CS equation and solving for $\beta_{\lambda}$ gives

$$
\begin{equation*}
\beta_{\lambda}=\frac{3 \lambda^{2}}{16 \pi^{2}}+O\left(\lambda^{3}\right) \tag{8.80}
\end{equation*}
$$

We observe that $\beta_{\lambda}$ at 1 -loop order is given by the coefficient of the $\frac{1}{\epsilon}$-singularity of $\delta_{\lambda}$ in dimensional regularization.

### 8.7 Computation of $\beta$-functions in massless theories

In this section we give general expressions for the beta-functions (8.69) in a massless scalar theory with dimensionless coupling ${ }^{4}$. Since by dimensional analysis, e.g. for the propagator,

$$
\begin{equation*}
i \mathcal{M}(p \rightarrow p)=a_{0} \Lambda^{2}+p^{2} \log \Lambda+\text { finite } \tag{8.81}
\end{equation*}
$$

the momentum independent quadratic divergence can be absorbed into the counter-term $\delta_{m}$. Thus in the sequel we only keep the logarithmic divergences and the counter-term $\delta_{Z}$.

- The 2-point function takes the form

$$
\begin{align*}
G_{2}(p) & =-\quad+1-\text { loop }+\longrightarrow{ }^{(1)}+\ldots \\
& =\frac{i}{p^{2}}+\frac{i}{p^{2}}\left(\frac{p^{2}}{i} A \log \frac{\Lambda^{2}}{-p^{2}}+\text { finite }\right) \frac{i}{p^{2}}+\frac{i}{p^{2}}\left(i p^{2} \delta_{Z}\right) \frac{i}{p^{2}}+\ldots \tag{8.82}
\end{align*}
$$

with

$$
\begin{equation*}
\delta_{Z}=A \log \frac{\Lambda^{2}}{\mu^{2}}+\text { finite } \tag{8.83}
\end{equation*}
$$

such that $G_{2}(p)$ is finite. We observe that

[^42]- the only explicit $\mu$ dependence is due to $\delta_{Z}$ at 1-loop, and that
$-\delta_{Z}$ is $O(\lambda)$.
As we will show below, $\beta(\lambda)=O\left(\lambda^{2}\right)$. Therefore at leading order the CS equation implies that

$$
\begin{equation*}
-\mu \frac{\partial}{\partial \mu} \delta_{Z}+2 \gamma_{\phi}=0 \tag{8.84}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\gamma_{\phi}=\frac{1}{2} \mu \frac{\partial}{\partial \mu} \delta_{Z}=-A . \tag{8.85}
\end{equation*}
$$

- More generally, the $n$-point function at 1 -loop order is of the form

$$
\begin{equation*}
G_{n}=\text { tree-level }+ \text { 1-loop PI }+ \text { vertex counter-term }+ \text { external legs, } \tag{8.86}
\end{equation*}
$$

where the external legs only contribute to $G$, not to the scattering amplitude. This can be parametrized as

$$
\begin{equation*}
G_{n}=\prod_{i} \frac{i}{p_{i}^{2}}\left[-i \lambda-i B \log \frac{\Lambda^{2}}{-p^{2}}-i \delta_{\lambda}+(-i \lambda) \sum_{i}\left(A_{i} \log \frac{\Lambda^{2}}{-p_{i}^{2}}-\delta_{Z_{i}}\right)\right], \tag{8.87}
\end{equation*}
$$

where $\sum_{i}$ extends over all $n$ external legs. To solve the CS equation

$$
\begin{equation*}
\left(\mu \frac{\partial}{\partial \mu}+\beta_{\lambda} \frac{\partial}{\partial \lambda}+\sum_{i=1}^{n} \gamma_{i}\right) G_{n}=0 \tag{8.88}
\end{equation*}
$$

we note that

- only $\delta_{\lambda}, \delta_{Z_{i}}$ contribute to $\mu \frac{\partial}{\partial \mu} G_{n}$ and that
- all terms except $-i \lambda$ are of order $\lambda^{2}$. So to leading order

$$
\begin{equation*}
\frac{\partial}{\partial \lambda} G_{n}=((-i)+O(\lambda)) \prod_{i} \frac{i}{p_{i}^{2}} . \tag{8.89}
\end{equation*}
$$

To leading order in $\lambda$ therefore

$$
\begin{equation*}
\mu \frac{\partial}{\partial \mu}\left(\delta_{\lambda}-\lambda \sum_{i} \delta_{Z_{i}}\right)+\beta(\lambda)+\lambda \sum_{i} \underbrace{\frac{1}{2} \mu \frac{\partial}{\partial \mu} \delta_{Z_{i}}}_{=\gamma_{i}}=0 . \tag{8.90}
\end{equation*}
$$

Since

$$
\begin{equation*}
\delta_{\lambda}=-B \log \frac{\Lambda^{2}}{\mu^{2}}+\text { finite, } \quad \delta_{Z_{i}}=A_{i} \log \frac{\Lambda^{2}}{\mu^{2}}+\text { finite } \tag{8.91}
\end{equation*}
$$

the $\beta$-function becomes

$$
\begin{equation*}
\beta(\lambda)=\mu \frac{\partial}{\partial \mu}\left(-\delta_{\lambda}+\frac{1}{2} \lambda \sum_{i} \delta_{Z_{i}}\right)=-2 B-\lambda \sum_{i} A_{i} . \tag{8.92}
\end{equation*}
$$

To leading order the finite terms do not enter the expressions for $\beta(\lambda)$ and $\gamma(\lambda)$. Since the different renormalisation schemes only differ by the finite terms, the 1-loop results are independent of the scheme. At next-to-leading order, the $\beta$-functions are scheme-dependent.

## 1-loop $\beta$-function of QED

As an application we compute the 1-loop $\beta$-function in QED. Recall our conventions for the counterterms appearing in the renormalisation of QED as discussed in QFT 1:

$$
\begin{align*}
& \delta_{3}=Z_{3}-1 \leftrightarrow \text { photon renormalisation }  \tag{8.93}\\
& \delta_{2}=Z_{2}-1 \leftrightarrow \text { fermion renormalisation }
\end{align*}
$$

and


To leading order in the coupling $e$, the $\gamma$-functions therefore are

$$
\begin{equation*}
\gamma_{3}=\frac{1}{2} \mu \frac{\partial}{\partial \mu} \delta_{3}, \quad \gamma_{2}=\frac{1}{2} \mu \frac{\partial}{\partial \mu} \delta_{2} \tag{8.94}
\end{equation*}
$$

and from the Callan-Symanzyk equation of the 3-point function one concludes that to leading order in $e$

$$
\begin{equation*}
\beta_{e}=\mu \frac{\partial}{\partial \mu}\left(-e \delta_{1}+\frac{1}{2} 2 e \delta_{2}+\frac{1}{2} e \delta_{3}\right) \tag{8.95}
\end{equation*}
$$

As will be discussed in more detail in the tutorials, the explicit expressions for $\delta_{1}, \delta_{2}, \delta_{3}$ obtained in QFT I ( with the renormalisation scale $\mu$ identified with the electron mass scale $m$ ) imply

$$
\begin{equation*}
\gamma_{2}=\gamma_{3}=\frac{e^{2}}{12 \pi^{2}}+O\left(e^{3}\right), \quad \beta=\frac{e^{3}}{12 \pi^{2}}+\mathcal{O}\left(e^{4}\right) \tag{8.96}
\end{equation*}
$$

in Feynman gauge.

### 8.8 The running coupling

Having computed the 1 -loop $\beta$-functions, we are in a position to discuss the physics encoded in these. Let us repeat our mantra that it is the renormalisation conditions which define the meaning of the physical coupling $\lambda$ at the renormalisation scale $\mu$. E.g. in massless $\phi^{4}$ theory we defined the dimensionless coupling $\lambda$ by declaring that

$$
\begin{equation*}
i \mathcal{M}\left(p_{1} p_{2} \rightarrow p_{3} p_{4}\right)=-i \lambda \quad \text { at } \quad s=t=u=-\mu^{2} \tag{8.97}
\end{equation*}
$$

The so-defined coupling is therefore really a function of the renormalisation scale $\mu, \lambda=\lambda(\mu)$, and this $\lambda(\mu)$ is the effective coupling relevant for processes with typical momenta of order $\mu$.
To appreciate this latter point, recall that $i \mathcal{M}\left(p_{1} p_{2} \rightarrow p_{3} p_{4}\right)$ in the above theory is

$$
\begin{equation*}
i \mathcal{M}=-i \lambda-i \frac{\lambda^{2}}{32 \pi^{2}}\left[\log \frac{-s}{\mu^{2}}+\log \frac{-t}{\mu^{2}}+\log \frac{-u}{\mu^{2}}\right]+O\left(\lambda^{3}\right) \tag{8.98}
\end{equation*}
$$

with $\lambda=\lambda(\mu)$ as in (8.97).

- For momenta $s, t, u \sim-\mu^{2}$, the loop-corrections are small, and the tree-level result $-i \lambda(\mu)$ is a good approximation. In this sense, $\lambda(\mu)$ is a measure for the strength of the interaction. This is what we call the effective coupling for processes at momentum scale $\mu$.
- At scales $|s|,|t|,|u| \gg \mu^{2}$ (or $\ll \mu^{2}$ ), the 1-loop correction becomes big due to the appearance of a so-called large log. We see that perturbation theory is an expansion not only in $\lambda$, but also in $\log \frac{-p^{2}}{\mu^{2}}$ with $p$ a typical momentum/energy scale of the process.

Thus, at scales $s, t, u \sim-\mu_{1}^{2}$, with $\mu_{1} \gg \mu$ or $\mu_{1} \ll \mu$, the coupling $\lambda(\mu)$ defined at scale $\mu$ is not what we would call the relevant coupling - it does not reflect the actual strength of the interaction because the tree-level and the 1 -loop contribution differ so much.

- At scales $s, t, u \sim-\mu_{1}^{2}$ with $\mu_{1} \gg \mu\left(\right.$ or $\left.\mu_{1} \ll \mu\right)$ a better expansion parameter for perturbation theory is $\lambda\left(\mu_{1}\right)$ defined via

$$
\begin{equation*}
\lambda\left(\mu_{1}\right)=-\left.\mathcal{M}\right|_{s=t=u=-\mu_{1}^{2}} . \tag{8.99}
\end{equation*}
$$

We therefore conclude that indeed
$\lambda(\mu)$ gives the strength of the interaction at energy scale $\mu$.
Now, the evolution of this coupling $\lambda(\mu)$ as we vary $\mu$ is precisely what is encoded in the $\beta$-function through

$$
\begin{equation*}
\mu \frac{\mathrm{d}}{\mathrm{~d} \mu} \lambda(\mu)=\frac{\mathrm{d}}{\mathrm{~d} \log \mu} \lambda(\mu)=\beta(\lambda) . \tag{8.100}
\end{equation*}
$$

The change in $\lambda(\mu)$ as we change $\mu$ is called renormalisation group (RG) flow for reasons that will become clear in section (8.10). The above RG flow equation describes the running of the coupling in the above sense.
Being a first-order differential equation, (8.100) uniquely specifies $\lambda(\mu)$ in terms of the value of $\lambda$ at a references scale $\mu^{*}$. Suppose therefore that $\lambda^{*}=\left.\lambda(\mu)\right|_{\mu=\mu^{*}}$ as determined by one measurement at $\mu^{*}$. Then (8.100) is solved by

$$
\begin{equation*}
\int_{\lambda^{*}}^{\lambda(\mu)} \frac{\mathrm{d} \lambda^{\prime}}{\beta\left(\lambda^{\prime}\right)}=\int_{\mu^{*}}^{\mu} \mathrm{d} \log \mu^{\prime}=\log \frac{\mu}{\mu^{*}} \tag{8.101}
\end{equation*}
$$

As an example consider the 1-loop $\beta$-function in QED,

$$
\begin{equation*}
\beta(e)=\frac{1}{12 \pi^{2}} e^{3} . \tag{8.102}
\end{equation*}
$$

The associated RG flow equation is easily integrated (see the tutorials) and gives for the fine-structure $\alpha \equiv \frac{e^{2}}{4 \pi}$

$$
\begin{equation*}
\alpha(\mu)=\frac{\alpha^{*}}{1-\frac{2}{3 \pi} \alpha^{*} \log \frac{\mu}{\mu^{*}}} \quad \text { with } \quad \alpha^{*}=\left.\alpha(\mu)\right|_{\mu=\mu^{*}} \tag{8.103}
\end{equation*}
$$

Observe that according to this analysis, $\alpha(\mu)$ increases as $\mu$ increases and even seems to diverge at $\mu=\mu^{*} \exp \left(\frac{3 \pi}{2 \alpha^{*}}\right)$.

## Alternative RG bevaviours

Depending on the sign of the $\beta$-function there are 3 qualitatively different types of behavior for the RG flow of the coupling as well as hybrids thereof. Let us first turn to the extreme cases:

1. If $\beta(\lambda)>0, \lambda(\mu)$ increases (decreases) as $\mu$ increases (decreases). If we start at a perturbative value $\lambda^{*}$ at $\mu^{*}$ and follow the RG flow for increasing $\mu$, then at some scale $\lambda(\mu)$ ceases to be perturbative and perturbation theory, in particular our perturbative evaluation of the $\beta$-function, is no more reliable. If by a non-perturbative analysis beyond that point one finds that $\beta(\lambda)>0$ $\forall \lambda$, then $\lambda(\mu)$ increases indefinitely. This can happen either asymptotically, i.e. $\lambda \rightarrow \infty$ as $\mu \rightarrow \infty$, or even for finite value of $\mu, \lambda \rightarrow \infty$ as $\mu \rightarrow \mu^{*}$. The latter instance is referred to as a Landau pole. In presence of a Landau pole, the theory is not well-defined at all energy scales. ${ }^{5}$ As we will understand in section 8.10 we should think of the theory as to be replaced by a different theory at some high energy scale which resolves the Landau-pole.

On the other hand, the theory is perturbatively well-defined in the IR, where $\lambda$ is small. For instance, if $\beta(\lambda) \rightarrow 0$ as $\lambda \rightarrow 0$ as would be the case e.g. for $\beta(\lambda)=c \lambda^{x}$ with $c>0$, then $\lambda(\mu) \rightarrow 0$ as $\mu \rightarrow 0$. The coupling thus approaches an IR fixed point, and the theory becomes free in the IR. The fixed point $\lambda=0$ is therefore called Gaussian fixed point.
Perturbatively $\beta(\lambda)>0$ for $\phi^{4}$ theory in $d=4$ and for $U(1)$ gauge theory in $d=4$, but the non-perturbative status is less clear. ${ }^{6}$
2. If $\beta(\lambda)<0, \lambda(\mu)$ decreases (increases) as $\mu$ increases (decreases). The theory is perturbative in the UV, but ceases to be perturbative in the IR.
If $\beta(\lambda) \rightarrow 0$ as $\lambda \rightarrow 0$ as would be the case e.g. for $\beta(\lambda)=c \lambda^{x}$ with $c<0$, we have that $\lambda \rightarrow 0$ as $\mu \rightarrow \infty$ in a controlled way (i.e. $\lambda$ flows to a Gaussian UV fixed-point) and the theory becomes free in the UV. This is called asymptotic freedom.
In $d=4$ the only known example for asymptotic freedom is Yang-Mills theory to be discussed in detail later.
3. If $\beta \equiv 0 \forall \mu, \lambda(\mu)$ is independent of $\mu$. The theory is thus scale-independent or conformal. Since the counter-terms do not induce any scale dependence, there cannot be any UV divergences altogether - the theory is finite. For example $\mathcal{N}=4$ Super-Yang-Mills theory is finite because all divergences cancel order by order in perturbation theory due to the large amount of supersymmetry.

Apart from these three extreme cases, it can happen that $\beta(\lambda)$ interpolates between 1., 2., 3., i.e. changes sign and thus hits zero at some value.

[^43]- If $\beta(\lambda)>0$ to lowest order in $\lambda$, but higher order terms, relevant for bigger $\lambda$ and thus (since we started with $\beta(\lambda)>0$ ) for larger $\mu$, have a negative sign, then $\beta(\lambda)$ becomes 0 at $\lambda=\lambda^{*}$ :


As $\mu \rightarrow \infty, \lambda \rightarrow \lambda^{*}$ with $\beta\left(\lambda^{*}\right)=0$. $\lambda^{*}$ is called non-trivial $\mathbf{U V}$ fixed point.

- Alternatively, $\beta(\lambda)$ could start negative for small values of $\lambda$ and hit 0 at $\lambda^{*}$ :


Then $\lambda \rightarrow \lambda^{*}$ as $\mu \rightarrow 0$ (non-trivial IR fixed point). Such behavior is known to occur e.g. in $\phi^{4}$ theory in $d<4$. As will be discussed in the tutorials,

$$
\begin{equation*}
\beta(\lambda)=(d-4) \lambda+\frac{3 \lambda^{2}}{16 \pi^{2}} \Rightarrow \beta\left(\lambda^{*}\right)=0 \quad \text { for } \quad \lambda^{*}=\frac{16 \pi^{2}}{3}(4-d), \tag{8.104}
\end{equation*}
$$

where $\lambda$ is here the suitably rescaled (dimensionless) quartic scalar coupling. Thus for $d<4$ the theory approaches the celebrated Wilson-Fisher fixed point $\lambda^{*}>0$ in the IR.

## Anomalous dimension and power-law behavior near fixed points

Near a non-trivial fixed point $\lambda^{*}$, the 2-point function of a scalar theory in $d=4$ behaves as

$$
\begin{equation*}
G_{2}(p)=C\left(\frac{1}{p^{2}}\right)^{1-\gamma_{\phi}\left(\lambda^{*}\right)} \tag{8.105}
\end{equation*}
$$

This can be proven by integrating the Callan-Symanzyk equation and using approximate scale - independence near the fixed point. For a derivation we refer to the tutorials. The propagator (8.105) looks like the propagator of the free theory,

$$
\begin{equation*}
\left.G_{2}(p)\right|_{\left.\right|_{\text {free }}} \sim \frac{1}{p^{2}}, \tag{8.106}
\end{equation*}
$$

except for the different power $\left(1-\gamma^{*}\right)$. Consequently $\gamma_{\phi}$ is called anomalous dimension of $\phi$ since near $\lambda^{*}$, the propagator $G_{2}(p)$ behaves like in a free theory with a field of mass dimension ( $1-$ $\gamma_{\phi}\left(\lambda^{*}\right)$ ). More generally, such power-law behaviour of the propagator is a sign of scale-invariance (as realised near $\lambda=\lambda^{*}$ where $\beta=0$ ).

### 8.9 RG flow of dimensionful operators

Our analysis of the running of the dimensionless coupling $\lambda$ can be generalised to dimensionful couplings. To this end suppose that we add to a renormalisable action (containing only dimensionless couplings) in $d$ dimensions a term

$$
\begin{equation*}
\int \mathrm{d}^{d} x C_{i} O_{d_{i}} \tag{8.107}
\end{equation*}
$$

where $O_{d_{i}}$ is some local operator ${ }^{7}$ of mass dimension $d_{i}$, and treat it as a perturbation. Then define $\gamma_{i}$ via the generalised Callan-Symanzyk equation:

$$
\begin{equation*}
0=\left[\mu \frac{\partial}{\partial \mu}+\beta \frac{\partial}{\partial \lambda}+n \gamma_{\phi}+\gamma_{i} C_{i} \frac{\partial}{\partial C_{i}}\right] G_{n}\left(x ; \lambda, C_{i} ; \mu\right) . \tag{8.108}
\end{equation*}
$$

Compared to $\beta \frac{\partial}{\partial \lambda}$ we have scaled out one power of $C_{i}$ in $\gamma_{i} C_{i} \frac{\partial}{\partial C_{i}}$ so that $\gamma_{i}$ is dimensionless. Clearly $\gamma_{i} \rightarrow 0$ as $\lambda \rightarrow 0$ and $C_{i} \rightarrow 0$.
Now, the coupling $C_{i}$ has mass dimension

$$
\begin{equation*}
\left[C_{i}\right]=d-d_{i}, \tag{8.109}
\end{equation*}
$$

so we can define a dimensionless coupling $g_{i}$ via

$$
\begin{equation*}
C_{i}=g_{i} \mu^{d-d_{i}} . \tag{8.110}
\end{equation*}
$$

[^44]Rewriting $G_{n}\left(x ; \lambda, C_{i} ; \mu\right)$ as $\tilde{G}_{n}\left(x ; \lambda, g_{i} ; \mu\right)$ introduces some extra explicit $\mu$ dependence, and the Callan-Symanczyk equation becomes (after omitting the tilde)

$$
\begin{equation*}
0=[\mu \frac{\partial}{\partial \mu}+\beta(\lambda) \frac{\partial}{\partial \lambda}+n \gamma_{\phi}+\underbrace{\left(\gamma_{i}+d_{i}-d\right) g_{i}}_{=: \beta_{i}} \frac{\partial}{\partial g_{i}}] G_{n}\left(x ; \lambda, g_{i} ; \mu\right), \tag{8.111}
\end{equation*}
$$

or, more symmetrically,

$$
\begin{equation*}
0=\left[\mu \frac{\partial}{\partial \mu}+\beta(\lambda) \frac{\partial}{\partial \lambda}+n \gamma_{\phi}+\beta_{i} \frac{\partial}{\partial g_{i}}\right] G_{n}\left(x ; \lambda, g_{i} ; \mu\right) . \tag{8.112}
\end{equation*}
$$

By definition, the dimensionless $g_{i}$ satisfies the RG flow equation

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} \log \mu} g_{i}=\beta_{i}=\left(\gamma_{i}+d_{i}-d\right) g_{i} . \tag{8.113}
\end{equation*}
$$

Suppose that $\lambda \ll 1$ and $g_{i} \ll 1$. Then necessarily $\gamma_{i} \ll 1$ and

$$
\begin{equation*}
\beta_{i} \rightarrow\left(d_{i}-d+\ldots\right) g_{i} \tag{8.114}
\end{equation*}
$$

so that (8.113) is solved to first approximation by

$$
\begin{equation*}
g_{i}(\mu)=g_{i}\left(\mu^{*}\right) e^{\left(d_{i}-d\right)\left(\ln \frac{\mu}{\mu^{*}}\right)}=g_{i}\left(\mu^{*}\right)\left(\frac{\mu}{\mu^{*}}\right)^{d_{i}-d} \tag{8.115}
\end{equation*}
$$

Note that if $d_{i}-d=0$, corresponding to a dimensionless coupling constant, this crude approximation gives no information about the behavior of the coupling and one has to compute the perturbative contributions to the beta-function even to determine the leading behaviour of the RG flow. We conclude:

- If $\left[C_{i}\right]=d-d_{i}>0$, i.e. $C_{i}$ is a super-renormalisable coupling, then

$$
\begin{equation*}
g_{i} \gg g_{i}\left(\mu^{*}\right) \quad \text { for } \quad \mu \ll \mu^{*} \tag{8.116}
\end{equation*}
$$

in a perturbative theory.

- If $\left[C_{i}\right]=d-d_{i}<0$, i.e. $C_{i}$ is a non-renormalisable coupling, then

$$
\begin{equation*}
g_{i}(\mu) \ll g_{i}\left(\mu^{*}\right) \tag{8.117}
\end{equation*}
$$

in a perturbative theory.
Non-renormalisable couplings become irrelevant in the IR: Even if they are present at higher scales, they become small for smaller energy processes. On the other hand, super-renormalisable couplings become relevant in the IR.

But it is important to keep in mind that if the coupling $\lambda$ is big, then $\gamma_{i}$ cannot be neglected in (8.113). In principle, this can change the qualitative behaviour of the flow.

For example, as already noted, in $\phi^{4}$ theory in $d$ dimensions the rescaled dimensionless coupling $\lambda$ behaves as

$$
\begin{equation*}
\beta(\lambda)=(d-4) \lambda+\frac{\lambda^{2}}{16 \pi^{2}} \tag{8.118}
\end{equation*}
$$

For $d<4 \lambda$ is super-renormalisable, but for big values of $\lambda$ due to the second term $\beta(\lambda)$ becomes positive and changes even the sign of the $\beta$-function.

### 8.10 Wilsonian effective action \& Renormalisation Semi-Group

Our philosophy behind renormalisation has so far been this:

- The cutoff $\Lambda$ is merely a way to regulate divergent integrals without any physical meaning.
- Renormalisation is a trick to remove the cutoff-dependence in physical amplitudes. This procedure allows us to take $\Lambda \rightarrow \infty$ without encountering any divergences.
- This comes at the cost of losing predictability for a number of physical masses and coupling.
- In a renormalisable theory, only a finite number of such physical couplings must be taken as input parameters from experiment, and we end up with a well-defined theory.

The so-called Wilsonian approach gives a rather different interpretation to renormalisation, inspired by analogies with condensed matter and statistical physics:

- According to Wilson, we should think of a QFT as an effective description of physics which is accurate only at energies below an intrinsic cutoff $\Lambda_{0}$ (or at distances larger than some $\Delta x_{\text {min }}$ ). At energies beyond $\Lambda_{0}$ the field theory picture does not correctly reflect the microscopic degrees of freedom.
- For instance, consider as our field the magnetisation density $s(\vec{x})$ of a metal in 3D statistical field theory. This $s(\vec{x})$ is the continuum limit of the spin degrees of freedom $s_{i}$ at lattice site $i$ of the atomic lattice. There is an intrinsic cutoff given by the lattice spacing $a$ :
At distances smaller than the lattice spacing $a$, no concept of continuous magnetisation exists - the continuous field theory description is strictly speaking invalid. The appearance of divergences in bare perturbation theory as we take $a \rightarrow 0$ reflects this mismatch of the field theory with the microscopic degrees of freedom.
Nonetheless, via renormalisation, one can take the continuum limit, but at the cost of having to take the couplings as input parameters. These parametrise our ignorance of the true physics at distances $\sim a$ from the perspective of the field theory.
- Applied to relativistic QFT, one would similarly argue that beyond some high energy scale $\Lambda_{0}$, QFT fails to capture the microscopic degrees of freedom. In particular, all matter gravitates,
and the effects of gravity become non-negligible at the (reduced) Planck scale

$$
\begin{equation*}
M_{\mathrm{pl}}=\frac{1}{\sqrt{8 \pi G_{N}}}=2.43 \times 10^{18} \mathrm{GeV} \tag{8.119}
\end{equation*}
$$

The observed UV divergences in bare perturbation theory result from the breakdown of our QFT beyond this scale. Still, in a renormalisable and asymptotically free (or at least UV safe) theory one can take the cutoff $\Lambda_{0} \rightarrow \infty$, but at the cost of parametrising our ignorance of the true degrees of freedom in a finite number of parameters.

- In the modern Wilsonian interpretation of Quantum Field Theory, we should view the QFT as an effective theory. E.g. in a non UV-finite QFT it is impossible to predict the vacuum energy $V_{0}$ - this is one of the input parameters. To explain the mystery of Dark Energy, we need a more fundamental, UV finite theory. To date, the only known theory that is UV finite and asymptotes to a weakly coupled effective Quantum Field Theory in the IR is String Theory. In String Theory, the concept of pointlike particles is abandoned in the UV: All particles are in fact excitations of a 1 -dimensional string of length $\ell_{s}$. At distances above $\ell_{s}$ we can approximate the theory by a point particle theory, i.e. a QFT as we have studied it. Thus the string length is the intrinsic cutoff of the low-energy effective QFT, and at distances smaller than $\ell_{s}$ the theory flows to a rather different (in particular non-local) description without UV divergences and without free dimensionless input parameters.

After this general outline we now demonstrate quantitatively how the high energy degrees of freedom are indeed captured in the effective coupling constants at lower energies. This will give a physical meaning to the running of the couplings. In view of the above motivation let us start with a bare Lagrangian in $d$ Euclidean dimensions, e.g.

$$
\begin{equation*}
\mathcal{L}=\frac{1}{2}(\partial \phi)^{2}+\frac{1}{2} m^{2} \phi^{2}+\frac{\lambda}{4!} \phi^{4} \tag{8.120}
\end{equation*}
$$

and impose a momentum cutoff $\Lambda_{0}$ : In path integral quantisation, this means we restrict the measure - written in Fourier space - to the Fourier modes $\tilde{\phi}(k)$ with $0 \leq|k| \leq \Lambda_{0}$, i.e. define the theory via

$$
\begin{equation*}
Z=\int[\mathcal{D} \phi]_{\Lambda_{0}} e^{-\int \mathrm{d}^{d} x_{E}(\mathcal{L}+J \cdot \phi)}, \tag{8.121}
\end{equation*}
$$

where $[\mathcal{D} \phi]_{\Lambda_{0}} \equiv \prod_{|k| \leq \Lambda_{0}} \mathrm{~d} \tilde{\phi}(k)$. Note the specific combinations of minus signs and absence of factors of $i$ in the Euclidean theory. This partition function enables us to compute amplitudes up to a typical energy scale $\Lambda_{0}$ of external particles. Suppose now that we are interested in processes at typical energy scales below

$$
\begin{equation*}
\Lambda=b \Lambda_{0} \quad 0<b<1 \tag{8.122}
\end{equation*}
$$

As we saw in the previous picture, loop computations will give rise to large log-corrections if $\Lambda \ll \Lambda_{0}$. To avoid these, we can integrate out the degrees of freedom between $\Lambda$ and $\Lambda_{0}$ once and for all and work with a path integral defined only in terms of the remaining degrees of freedom below $\Lambda$. To compute this new effective path integral we first split each Fourier mode as

$$
\begin{equation*}
\tilde{\phi}(k)=\phi(k)+\hat{\phi}(k), \tag{8.123}
\end{equation*}
$$

where

$$
\hat{\phi}(k)= \begin{cases}0 & \text { if } 0 \leq|k| \leq b \Lambda_{0}  \tag{8.124}\\ \tilde{\phi}(k) & \text { if } b \Lambda_{0} \leq|k| \leq \Lambda_{0}\end{cases}
$$

and

$$
\phi(k)= \begin{cases}\tilde{\phi}(k) & \text { if } 0 \leq|k| \leq b \Lambda_{0}  \tag{8.125}\\ 0 & \text { if } b \Lambda_{0} \leq|k| \leq \Lambda_{0}\end{cases}
$$

Taking $J=0$ for brevity, we get

$$
\begin{align*}
Z= & \int \mathcal{D} \phi \mathcal{D} \hat{\phi} e^{-\int \mathrm{d}^{d} x\left[\frac{1}{2}(\partial \phi+\partial \hat{\phi})^{2}+\frac{1}{2} m^{2}(\phi+\hat{\phi})^{2}+\frac{\lambda}{4!}(\phi+\hat{\phi})^{4}\right]} \\
= & \int \mathcal{D} \phi e^{-\int \mathrm{d}^{d} x\left(\frac{1}{2} \partial \phi^{2}+\frac{1}{2} m^{2} \phi^{2}+\frac{\lambda}{4} \phi^{4}\right)}  \tag{8.126}\\
& \times \int \mathcal{D} \hat{\phi} e^{-\int \mathrm{d}^{d} x\left[\frac{1}{2} \partial \hat{\phi}^{2}+\frac{1}{2} m^{2} \hat{\phi}^{2}+\frac{\lambda}{4!} \hat{\phi}^{4}+\frac{\lambda}{6} \phi \hat{\phi}^{3}+\frac{\lambda}{4} \phi^{2} \hat{\phi}^{2}+\frac{\lambda}{6} \hat{\phi} \phi^{3}\right]},
\end{align*}
$$

where terms $\phi \hat{\phi}$ are omitted due to orthogonality of different momentum modes. The partition function is of the form

$$
\begin{equation*}
Z=\int[\mathcal{D} \phi]_{b \Lambda_{0}} e^{-S[\phi]} \int[\mathcal{D} \hat{\phi}]_{b \Lambda_{0} \leq|k| \leq \Lambda_{0}} e^{-\hat{S}[\phi, \hat{\phi}]} . \tag{8.127}
\end{equation*}
$$

The key idea is that this can be rewritten as

$$
\begin{equation*}
Z=\int[\mathcal{D} \phi]_{b \Lambda_{0}} e^{-\int \mathrm{d}^{d} x \mathcal{L}_{\text {eff }}} \tag{8.128}
\end{equation*}
$$

To this end we view all terms in $\hat{S}[\phi, \hat{\phi}]$ except the kinetic term as a perturbation, which is justified for

$$
\begin{equation*}
\lambda \ll 1 \text { and } m \ll \Lambda=b \Lambda_{0} . \tag{8.129}
\end{equation*}
$$

Then perturbatively

$$
\begin{align*}
\int \mathcal{D} \hat{\phi} e^{-\hat{S}[\phi, \hat{\phi}]} & =\int \mathcal{D} \hat{\phi} e^{-\int \mathrm{d}^{d} x \frac{1}{2} \partial \hat{\phi}^{2}}\left[1-\int \mathrm{d}^{d} x\left(\frac{1}{2} m^{2} \hat{\phi}^{2}+\frac{\lambda}{4!} \hat{\phi}^{4}+\frac{\lambda}{6} \phi \hat{\phi}^{3}+\frac{\lambda}{4} \phi^{2} \hat{\phi}^{2}+\frac{\lambda}{6} \hat{\phi} \phi^{3}\right)+\ldots\right] \\
& \equiv\langle[\ldots]\rangle . \tag{8.130}
\end{align*}
$$

This is computed via Wick's theorem in perturbation theory, where we contract only pairs of $\hat{\phi}$, and in the resulting propagator only momenta $b \Lambda_{0} \leq|k| \leq \Lambda$ appear. Let us sketch how this works:

1. Consider for example the term

$$
\begin{align*}
&\left\langle-\int \mathrm{d}^{d} x \frac{\lambda}{4} \phi^{2} \hat{\phi}^{2}\right\rangle=\int \mathcal{D} \hat{\phi} e^{-\int \mathrm{d}^{d} x \frac{1}{2} \hat{\phi} \hat{\phi}^{2}}\left(-\int \mathrm{d}^{d} x \frac{\lambda}{4} \phi^{2}(x) \hat{\phi}^{2}(x)\right) \\
&=-\int \mathrm{d}^{d} x \frac{1}{2} \phi^{2}(x) \frac{\lambda}{2} \underbrace{\hat{\phi}(x) \hat{\phi}(x)}  \tag{8.131}\\
&=\int \frac{d^{d} p}{\left(2 \pi \pi d^{d}\right.} e^{-i p(x-x)} \\
&=\int \mathrm{d}^{d} x\left(-\left.\frac{1}{p^{2}}\right|_{b \Lambda_{0} \leq|p| \leq \Lambda_{0}} \phi^{2}(x) \mu\right)
\end{align*}
$$

with

$$
\begin{equation*}
\mu=\left.\frac{\lambda}{2} \int \frac{\mathrm{~d}^{d} p}{(2 \pi)^{d}} \frac{1}{p^{2}}\right|_{b \Lambda_{0} \leq|p| \leq \Lambda_{0}}=\frac{\lambda}{(4 \pi)^{d / 2} \Gamma(d / 2)} \frac{1-b^{d-2}}{d-2} \Lambda^{d-2} . \tag{8.132}
\end{equation*}
$$

Note that in the Euclidean theory no factor of $i$ appears in the propagator. Graphically one represents the $\hat{\phi}$-propagator with two extra lines:

$$
\hat{\phi}(x) \hat{\phi}(y) \equiv x \longmapsto
$$



Together with $n$-fold powers of that same term (which arise from expanding the exponential in $\left.e^{-[S(\phi, \hat{\phi})]}\right)$ this just gives a correction to the mass term present already in $e^{-S[\phi]}$.
2. Next consider the square of this term:

$$
\frac{1}{4}\left\langle\frac{\lambda}{4} \phi^{2} \hat{\phi}^{2} \frac{\lambda}{4} \phi^{2} \hat{\phi}^{2}\right\rangle \sim(
$$



The first diagram is just the square of the mass correction, as advocated. The second term is schematically of the form

$$
\begin{equation*}
\int d^{d} x d^{d} y \phi^{2}(x) \phi^{2}(y) \simeq \int d^{d} x\left(\phi^{4}(x)+c \phi^{2}(x)(\partial \phi(x))^{2}+\ldots\right. \tag{8.133}
\end{equation*}
$$

The first summand provides a contribution to the $\phi^{4}$ coupling in $S[\phi]$, while the second gives a new, derivative interaction not present in our original action.
3. Likewise, also higher order interaction terms in $\phi$ appear, e.g. the contraction

provides a $\phi^{6}$ interaction in $S_{\text {eff }}$.

This way, all possible interactions are created in $S_{\text {eff }}$ which are consistent with the symmetries of the bare Lagrangian (here just $\phi \rightarrow-\phi$ ), and we find

$$
\begin{equation*}
Z=\int[\mathcal{D} \phi]_{b \Lambda_{0}} e^{-S_{\mathrm{eff}}^{b \Lambda_{0}}} \tag{8.134}
\end{equation*}
$$

with

$$
\begin{gather*}
\int \mathrm{d}^{d} x \mathcal{L}_{\text {eff }}^{b \Lambda_{0}}=\int \mathrm{d}^{d} x\left[\frac{1}{2}(1+\Delta Z)(\partial \phi)^{2}+\frac{1}{2}\left(m^{2}+\Delta m^{2}\right) \phi^{2}+\frac{1}{4!}(\lambda+\Delta \lambda) \phi^{4}\right.  \tag{8.135}\\
\left.+\Delta C\left(\partial_{\mu} \phi\right)^{4}+\Delta D \phi^{6}+\ldots\right] \tag{8.136}
\end{gather*}
$$

The object $S_{\text {eff }}^{b \Lambda_{0}}$ is called Wilsonian effective action. The coefficients $\Delta Z, \Delta m, \Delta \lambda, \Delta C, \Delta D$ can be systematically computed in perturbation theory. They receive contributions only from the momentum regime $b \Lambda_{0} \leq|k| \leq \Lambda_{0}$. When we compute correlators at scales below $\Lambda=b \Lambda_{0}$ via the Wilsonian effective action (8.134), only momenta $|k| \leq \Lambda=b \Lambda_{0}$ appear in the loops since all effects of the modes with $b \Lambda_{0}<|k|<\Lambda_{0}$ are already encoded in $\mathcal{L}_{\text {eff }}$.
The Wilsonian effective action is not to be confused with the quantum effective action $\Gamma[\varphi]$. As discussed in section 7.9, the latter includes all quantum effects, i.e. computation with $\Gamma[\varphi]$ at treelevel (no loops!) gives the full quantum theory. By contrast, in the Wilsonian effective action $S_{\text {eff }}^{b \Lambda_{0}}$ only the quantum effects due to modes from $b \Lambda_{0}<|k|<\Lambda_{0}$ are included, and in computations with $S_{\text {eff }}^{b \Lambda_{0}}$ loops with momenta $0 \leq|k| \leq \Lambda=b \Lambda_{0}$ must still be performed.
For better comparison of $\mathcal{L}_{\text {eff }}^{b \Lambda_{0}}$ with the original Lagrangian, which we denote now by $\mathcal{L}^{\Lambda_{0}}$, we bring the kinetic term to standard form by rescaling

$$
\begin{equation*}
\phi^{\prime}=\phi(1+\Delta Z)^{1 / 2} . \tag{8.137}
\end{equation*}
$$

Then the partition function $Z$ is given by

$$
\begin{equation*}
Z=\int\left[\mathcal{D} \phi^{\prime}\right]_{\Lambda} e^{-S_{\mathrm{eff}}^{\prime}} \quad \text { with } \quad \Lambda=b \Lambda_{0} \tag{8.138}
\end{equation*}
$$

and

$$
\begin{equation*}
S_{\mathrm{eff}}^{\prime \Lambda}=\int \mathrm{d}^{d} x\left[\frac{1}{2}\left(\partial \phi^{\prime}\right)^{2}+\frac{1}{2} m^{\prime 2} \phi^{\prime 2}+\frac{\lambda^{\prime}}{4!} \phi^{\prime 4}+C^{\prime}\left(\partial \phi^{\prime}\right)^{4}+D^{\prime} \phi^{\prime 6}+\ldots\right] \tag{8.139}
\end{equation*}
$$

for

$$
\begin{align*}
m^{\prime 2} & =\left(m^{2}+\Delta m^{2}\right)(1+\Delta Z)^{-1}, \\
\lambda^{\prime} & =(\lambda+\Delta \lambda)(1+\Delta Z)^{-2}, \\
C^{\prime} & =(C+\Delta C)(1+\Delta Z)^{-2},  \tag{8.140}\\
D^{\prime} & =(D+\Delta D)(1+\Delta Z)^{-3} .
\end{align*}
$$

Note that in our case $C=D=0$ in the original $\mathcal{L}^{\Lambda_{0}}$. Thus the effect of integrating out the degrees of freedom with $b \Lambda_{0}<|k|<\Lambda_{0}$ amounts to renormalising the couplings in the Lagrangian as above.

The successive application of this operation gives rise to the renormalisation group. This is really only a semi-group, since this procedure only allows us to lower the cutoff by integrating out degrees of freedom, i.e. there does not exist an inverse element.

The running of the couplings in the Wilsonian picture is to be interpreted as the evolution of the cutoff dependent couplings in the effective action as we change the cutoff. To make a connection to our previous picture we must identify the cutoff $\Lambda$ with the renormalisation scale $\mu$.
I.e. the effective couplings $\lambda(\Lambda)$ are defined by specifying observables computed from the effective action with cutoff $\Lambda$. This replaces our previous renormalisation condition fixing $\lambda(\mu)$.
With this identification, the fundamental form of the running couplings agree to leading order in perturbation. This is all we can hope for, because, as always, beyond leading order the specific scale dependence of the coupling depends on the specific renormalisation scheme. E.g. for $\phi^{4}$ theory in $d=4$ carrying out the Wilsonian renormalisation procedure as sketched above gives

$$
\begin{equation*}
\left.\Delta Z\right|_{1-\text { loop }}=0,\left.\quad \Delta \lambda\right|_{1-\text { loop }}=-\frac{3 \lambda^{2}}{16 \pi^{2}} \log \frac{1}{b} \tag{8.141}
\end{equation*}
$$

With $b=\Lambda / \Lambda_{0}<1$ this gives for small $|1-b|$

$$
\begin{equation*}
\lambda(\Lambda)=\lambda\left(\Lambda_{0}\right)+\frac{3}{16 \pi^{2}} \lambda^{2}\left(\Lambda_{0}\right) \log \frac{\Lambda}{\Lambda_{0}} \tag{8.142}
\end{equation*}
$$

in agreement with our previous results for the $\beta$-function of $\phi^{4}$-theory. More generally, for a dimensionful coupling $\int \mathrm{d}^{d} x C_{d_{i}} O_{d_{i}}$ with

- $O_{d_{i}}$ a local operator of mass dimension $d_{i}$ and
- $C_{d_{i}}$ of mass dimension $\left[C_{d_{i}}\right]=d-d_{i}$
one finds in $\phi^{4}$-theory

$$
\begin{equation*}
C_{d_{i}}(\Lambda)=C_{d_{i}}\left(\Lambda_{0}\right)+(\text { number }) \lambda^{d_{i} / 2}\left(\Lambda_{0}\right)\left(\Lambda^{d-d_{i}}-\Lambda_{0}^{d-d_{i}}\right)+\text { higher order. } \tag{8.143}
\end{equation*}
$$

In fact, this is just as expected by dimensional analysis. As in our previous approach it is useful to define a dimensionless coupling, but now in terms of $\Lambda$,

$$
\begin{equation*}
g_{i}(\Lambda)=C_{i}(\Lambda) \Lambda^{d_{i}-d} \tag{8.144}
\end{equation*}
$$

which enjoys the RG flow

$$
\begin{equation*}
g_{i}(\Lambda)=\left[g_{i}\left(\Lambda_{0}\right)+(\text { number }) \lambda^{d_{i} / 2}\left(\Lambda_{0}\right)\left(\left(\frac{\Lambda_{0}}{\Lambda}\right)^{d_{i}-d}-1\right)+\ldots\right]\left(\frac{\Lambda_{0}}{\Lambda}\right)^{d-d_{i}} \tag{8.145}
\end{equation*}
$$

Provided the theory is perturbative, the scaling due to

$$
\begin{equation*}
\left(\frac{\Lambda_{0}}{\Lambda}\right)^{d-d_{i}}=\left(\frac{\Lambda_{0}}{\Lambda}\right)^{\left[C_{i}\right]} \tag{8.146}
\end{equation*}
$$

dominates and thus if $C_{i}$ is non-renormalisable $\left(\left[C_{i}\right]<0\right)$, the initial value of $g_{i}\left(\Lambda_{0}\right)$ is irrelevant at scales $\Lambda \ll \Lambda_{0}$ and

$$
\begin{equation*}
C_{i}(\Lambda)=O\left(\lambda^{d_{i} / 2}\right) \Lambda^{\left[C_{i}\right]} \ll \Lambda^{\left[C_{i}\right]} \tag{8.147}
\end{equation*}
$$

for $\lambda \ll 1$. This agrees with our conclusions from the Callan-Symanzik approach: Non-renormalisable couplings become irrelevant in the IR: Even if they are present at higher scales, they become small for smaller energy processes. On the other hand, super-renormalisable couplings become relevant in the IR.

## Conclusion

- In the modern Wilsonian interpretation, a QFT is defined via an intrinsic cutoff scale $\Lambda_{0}$ beyond which it does not capture the microscopic degrees of freedom any more. For particle physics one would assume that at the very least

$$
\begin{equation*}
\Lambda_{0} \sim M_{\mathrm{pl}} \sim 2.43 \times 10^{18} \mathrm{GeV} \tag{8.148}
\end{equation*}
$$

which is the scale where quantum gravity effects become non-negligibly strong such that it is no longer justified to ignore them.

- Let us define our theory at this intrinsic cutoff by a Lagrangian $\mathcal{L}^{\Lambda_{0}}$. We can think of the couplings $C_{i}\left(\Lambda_{0}\right)$ as being computed by whatever more fundamental theory takes over at scales above $\Lambda_{0}$. In the spirit of an effective description, we view these $C_{i}\left(\Lambda_{0}\right)$ as parametrising the physics beyond $\Lambda_{0}$. At scales $\Lambda_{0}$, all possible non-renormalisable couplings are important, because the underlying microscopic degrees of freedom become important at $\Lambda_{0}$ and it takes a lot of parameters to describe these in QFT. I.e. all couplings consistent with the symmetries of the theory will appear in $\mathcal{L}^{\Lambda_{0}}$.
- For processes at our scales $\Lambda \ll \Lambda_{0}$ we compute with the effective action $\mathcal{L}_{\mathrm{eff}}^{\Lambda}$ with couplings $C_{i}(\Lambda)$ arising by RG flow starting from $C_{i}\left(\Lambda_{0}\right)$. Crucially, as long as the theory is weakly coupled up to $\Lambda_{0}$, the non-renormalisable couplings at $\Lambda \ll \Lambda_{0}$ are of order

$$
\begin{equation*}
C_{i}(\Lambda) \sim \lambda^{d_{i} / 2} \Lambda^{\left[C_{i}\right]} \ll \Lambda^{\left[C_{i}\right]} \tag{8.149}
\end{equation*}
$$

whatever they are originally at $\Lambda_{0}$. Since the microscopic degrees of freedom become important only at high scales, it takes only a few parameters - the renormalisable and super-renormalisable ones - to effectively approximate the physics at $\Lambda \ll \Lambda_{0}$. This explains why we only observe renormalisable couplings at scale $\Lambda$.

- Despite the physical meaning of $\Lambda_{0}$ (the intrinsic cutoff) we can formally take $\Lambda_{0} \rightarrow \infty$ provided all couplings remain finite, i.e. no Landau poles arise. This was our original motivation in doing renormalisation, but from a Wilsonian perspective we now understand that our theory remains an effective theory to the extent that all couplings are really input parameters and cannot be computed from first principles without knowing the underlying theory.


## On the fine-tuning problem in scalar theories

For scalar theories, in particular $\phi^{4}$ in $d=4$, there remains a problem because of the RG flow of the scalar mass:

$$
\begin{equation*}
m^{2}(\Lambda)=m^{2}\left(\Lambda_{0}\right)+(\text { number }) \lambda\left(\Lambda_{0}\right)\left(\Lambda^{2}-\Lambda_{0}^{2}\right)+\ldots, \tag{8.150}
\end{equation*}
$$

which is the special case of (8.145) for $d-d_{i}=2$. At scales $\Lambda \ll \Lambda_{0}$, the dominant term is given by $\lambda \Lambda_{0}^{2}$, and thus $m^{2}(\Lambda)$ is of order $\lambda\left(\Lambda_{0}\right) \Lambda_{0}^{2}$. This is much bigger than we naturally expect in an effective theory at scale $\Lambda$. To arrive at a value for $m^{2}(\Lambda)$ of order $\Lambda^{2}$ instead, it would appear one must carefully fine-tune the tree-level $m^{2}\left(\Lambda_{0}\right)$ such as to almost precisely cancel the loop corrections of order $\Lambda_{0}^{2}$.
For example, in the Standard Model of Particle Physics, the Higgs particle is a scalar field with $m_{H} \sim 126 \mathrm{GeV}$. But if we take $\Lambda_{0} \sim 10^{19} \mathrm{GeV}$, then

$$
\begin{equation*}
\frac{m_{H}^{2}}{\Lambda_{0}^{2}} \cong 10^{-32} \tag{8.151}
\end{equation*}
$$

The necessary fine-tuning to achieve this is considered 'technically unnatural' and gives rise to the so-called hierarchy problem: What stabilises the mass of a scalar at hierachically smaller mass than $M_{\mathrm{pl}}$ ?
Basically there are 3 alternative types of solutions to the fine-tuning problem that are being discussed in the modern literature:

1. If the inherent cutoff for the Standard Model were not $M_{P l}$, but at, say, 1 TeV , the fine-tuning would be dramatically reduced. This would imply that new physics takes over at $\Lambda_{0} \sim 1 \mathrm{TeV}$. Candidates for such new dynamics are supersymmetry (SUSY) or large extra dimensions. This reasoning motivates the search for such new physics at LHC.
2. Alternatively we could consider accepting the fine-tuning, which is not logically impossible for instance, but not necessarily, motivated by the so-called anthropic principle (which would state that only values of $m_{H} \sim 1 \mathrm{TeV}$ could give rise to structure formation as we observe it and thus, by definition, only a universe in which $m_{H}$ is so fine-tuned can be observed). Clearly this is hard to make precise, but it is a logical possibility in principle.
3. Yet another approach would be to completely revise our notion of naturalness, even though it is unclear at present what it should be replaced by.

It is important to appreciate that such a hierarchy problem does not arise for fermions and vectors because their masses are stabilised by extra symmetry, namely gauge symmetry for vectors and chiral symmetry for fermions. Small masses for these fields are technically natural in the sense that once we choose them small at tree-level, they remain small including loop-corrections. This is guaranteed whenever setting the mass to zero implies an enhancement of symmetry - such as in the present case gauge symmetry for massless vector bosons and chiral symmetry for massless fermions.

## Chapter 9

## Quantisation of Yang-Mills-Theory

### 9.1 Recap of classical YM-Theory

Let us quickly recall the description of classical non-abelian gauge theory as introduced in QFT1. Our starting point is a Lie group $H$ with Lie algebra Lie $(H)$ such that every element $h \in H$ can be expressed as

$$
\begin{equation*}
h=e^{-i g \alpha} \quad \text { with } \quad \alpha \in \operatorname{Lie}(H), g \in \mathbb{R} \tag{9.1}
\end{equation*}
$$

In terms of a basis $T^{a}$ of the Lie algebra, i.e. in terms of a complete set of generators of the Lie group $H, \alpha$ can be written as

$$
\begin{equation*}
\alpha=\sum_{a} \alpha_{a} T^{a} \equiv \alpha_{a} T^{a} \tag{9.2}
\end{equation*}
$$

The $T^{a}$ satisfy the defining relations

$$
\begin{equation*}
\left[T^{a}, T^{b}\right]=i f_{c}^{a b} T^{c}=-\left[T^{b}, T^{a}\right] \tag{9.3}
\end{equation*}
$$

which fulfill the Jacobi identity

$$
\begin{equation*}
\left[\left[T^{a}, T^{b}\right], T^{c}\right]+\left[\left[T^{c}, T^{a}\right], T^{b}\right]+\left[\left[T^{b}, T^{c}\right], T^{a}\right]=0 \tag{9.4}
\end{equation*}
$$

A Lie algebra possesses a symmetric bilinear from, the so-called Killing form,

$$
\begin{equation*}
\kappa^{a b}=T^{a} \circ T^{b}=T^{b} \circ T^{a} \in \mathbb{R} \tag{9.5}
\end{equation*}
$$

which is invariant under the adjoint action of the Lie group $H$ on its Lie algebra, i.e.

$$
\begin{equation*}
h T^{a} h^{-1} \circ h T^{b} h^{-1}=T^{a} \circ T^{b} \forall h \tag{9.6}
\end{equation*}
$$

For an abstractly given Lie algebra, the Killing form is defined as

$$
\begin{equation*}
\kappa^{a b}=\operatorname{tr}\left[T^{a},\left[T^{b}, \cdot\right]\right] \equiv \sum_{c}\left(T_{c}^{*}\right)\left[T^{a},\left[T^{b}, T^{c}\right]\right] \tag{9.7}
\end{equation*}
$$

where $\left(T_{c}^{*}\right)$ is an element of a basis of the dual vector space such that $\left(T_{c}^{*}\right) T^{d}=\delta_{c}^{d}$.

Expanding $T^{a} \circ T^{b}=h T^{a} h^{-1} \circ h T^{b} h^{-1}$ with $h=e^{-i g T_{c}}$ to order $g$ gives

$$
\begin{equation*}
\left[T^{a}, T^{c}\right] \circ T^{b}=-T^{a} \circ\left[T^{c}, T^{b}\right] \tag{9.8}
\end{equation*}
$$

If $H$ is compact as a manifold, then the invariant from $\circ$ can be shown to be positive definite, and one can normalise it as

$$
\begin{equation*}
T^{a} \circ T^{b}=\frac{1}{2} \delta^{a b} \tag{9.9}
\end{equation*}
$$

The Killing form $\kappa^{a b}=T^{a} \circ T^{b}$ and its inverse can be used to raise and lower Lie-algebra indices. Having said this, we will from now on put all Lie algebra indices upstairs with summation over repeated indices understood. Then the structure constants $f^{a b c}$ are totally antisymmetric (as will be shown in the exercises) and invariant under cyclic permutations.
In the sequel we will typically think of $T^{a}$ as given in the defining matrix representation, i.e. as an $N \times N$ matrix acting on $\mathbb{C}^{N}$ for some $N$. E.g. for $H=\mathrm{SU}(k), N=k$ and the $T^{a}$ are represented as hermitian traceless $k \times k$ matrices. In such a matrix representation we can identify

$$
\begin{equation*}
T^{a} \circ T^{b} \equiv \operatorname{tr} T^{a} T^{b} \tag{9.10}
\end{equation*}
$$

Consider now a $\operatorname{Lie}(H)$-valued gauge potential, i.e. a connection

$$
\begin{equation*}
A_{\mu}(x) \equiv A_{\mu}^{a}(x) T^{a} \tag{9.11}
\end{equation*}
$$

transforming under a gauge transformation

$$
\begin{equation*}
U(x)=e^{-i g \alpha(x)} \tag{9.12}
\end{equation*}
$$

as

$$
\begin{equation*}
A_{\mu}(x) \mapsto U A_{\mu} U^{-1}+\frac{i}{g}\left(\partial_{\mu} U\right) U^{-1} \tag{9.13}
\end{equation*}
$$

To linear order this is

$$
\begin{equation*}
A_{\mu}(x) \mapsto A_{\mu}+D_{\mu} \alpha(x)+O\left(\alpha^{2}\right) \tag{9.14}
\end{equation*}
$$

where the adjoint covariant derivative acts on $\operatorname{Lie}(H)$-valued fields as

$$
\begin{equation*}
D_{\mu} \alpha(x):=\partial_{\mu} \alpha(x)+i g\left[A_{\mu}(x), \alpha(x)\right] \tag{9.15}
\end{equation*}
$$

In components this means that

$$
\begin{equation*}
A_{\mu}^{a}(x) \mapsto A_{\mu}^{a}(x)+\partial_{\mu} \alpha^{a}(x)+g f^{a b c} \alpha^{b}(x) A_{\mu}^{c}(x) \tag{9.16}
\end{equation*}
$$

The curvature associated with the connection $A_{\mu}(x)$ is

$$
\begin{equation*}
F_{\mu \nu}(x)=\frac{1}{i g}\left[D_{\mu}, D_{\nu}\right] \equiv F_{\mu \nu}^{a} T^{a} \tag{9.17}
\end{equation*}
$$

where we think of $\frac{1}{i g}\left[D_{\mu}, D_{\nu}\right]$ as acting on a $\operatorname{Lie}(H)$-valued field. Plugging in the above definition of the adjoint covariant derivative yields

$$
\begin{equation*}
F_{\mu \nu}(x)=\partial_{\mu} A_{v}-\partial_{v} A_{\mu}+i g\left[A_{\mu}, A_{v}\right] \tag{9.18}
\end{equation*}
$$

or in components

$$
\begin{equation*}
F_{\mu \nu}^{a}(x)=\partial_{\mu} A_{\nu}^{a}-\partial_{\nu} A_{\mu}^{a}-g f^{a b c} A_{\mu}^{b} A_{\nu}^{c} \tag{9.19}
\end{equation*}
$$

The field strength transforms in the adjoint representation of $H$ as

$$
\begin{equation*}
F_{\mu \nu} \mapsto U F_{\mu \nu} U^{-1} \tag{9.20}
\end{equation*}
$$

and satisfies the Bianchi identity

$$
\begin{equation*}
D_{\mu} F_{\alpha \beta}+D_{\beta} F_{\mu \alpha}+D_{\alpha} F_{\beta \mu}=0 \tag{9.21}
\end{equation*}
$$

The pure Yang-Mills Lagrangian

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{2} \operatorname{tr} F_{\mu \nu} F^{\mu \nu}=-\frac{1}{4} \sum_{a} F_{\mu \nu}^{a} F^{\mu v a} \tag{9.22}
\end{equation*}
$$

is gauge invariant. The equations of motion for $A_{\mu}$ read

$$
\begin{equation*}
D_{\mu} F^{\mu v}=0 \equiv \partial_{\mu} F^{\mu v}+i g\left[A_{\mu}, F^{\mu v}\right] \tag{9.23}
\end{equation*}
$$

The commutator introduces cubic and quartic interactions which will be discussed in detail later on.

### 9.2 Gauge fixing the path integral

Gauge invariance complicates the quantisation of Yang-Mills theory:

- In canonical quantisation, we would start by defining the conjugate momenta to $A_{\mu}$, but while

$$
\begin{equation*}
\Pi^{i}=\frac{\partial \mathcal{L}}{\partial \dot{A}_{i}} \sim F^{0 i} \tag{9.24}
\end{equation*}
$$

one finds

$$
\begin{equation*}
\Pi^{0}=\frac{\partial \mathcal{L}}{\partial \dot{A}_{0}}=0 \tag{9.25}
\end{equation*}
$$

This is a result of the fact that $A_{0}$ appears without a time derivative in $\mathcal{L}$ - it is a non-dynamical field. Variation of the action with respect to $A_{0}$ gives

$$
\begin{equation*}
D_{i} F^{0 i}=0 \tag{9.26}
\end{equation*}
$$

which is a non-dynamical constraint, and $A_{0}$ is a Lagrange multiplier enforcing this constraint. Thus canonical quantisation requires the technology ${ }^{1}$ of quantisation of constrained, or singular, systems. Quantisation of singular systems is a rich subject by itself. A variant of canonical

[^45]quantization with constraints for the simple case of $U(1)$ theory is in fact the Gupta-Bleuler quantisation method discussed in QFT1, but for non-abelian gauge theories things are more complicated due to the more involved form of the constraints. We will not start with canonical quantization, but rather begin with path-integral quantization and return to the problem of canonical quantization afterwards.

- Alternatively in path integral quantisation we would write the action as

$$
\begin{equation*}
S=\int \mathrm{d}^{4} x\left(-\frac{1}{2} A^{a v}\left(K \cdot A^{a}\right)_{v}+O\left(A^{3}\right)\right) \tag{9.27}
\end{equation*}
$$

where $(K \cdot A)_{v}=-\partial^{2} A_{v}+\partial_{\nu} \partial^{\mu} A_{\mu}$, find the propagator $i D_{F}=K^{-1}$ and perturbatively quantise the interacting theory. However, the operator $K$ is not invertible since it has a non-trivial kernel:

$$
\begin{equation*}
(K \cdot \partial \alpha)_{v}=0 \quad \forall \alpha(x) . \tag{9.28}
\end{equation*}
$$

In fact, this is a consequence of gauge invariance. Recall that an infinitesimal gauge transformation acts as

$$
\begin{equation*}
A_{\mu} \mapsto A_{\mu}+\partial_{\mu} \alpha+i g\left[A_{\mu}, \alpha\right] . \tag{9.29}
\end{equation*}
$$

Since the coupling $g$ enters only the interactions, the transformation $A_{\mu} \mapsto A_{\mu}+\partial_{\mu} \alpha+O(g)$ must leave the kinetic term invariant, which is nothing other than the statement (9.28). Note that this problem is therefore independent of whether we consider abelian or non-abelian gauge theory.

We will now begin quantization of Yang-Mills theory in the path-integral approach. Apart from providing a general method to compute correlation functions of gauge invariant operators, this will result in a formalism in which we can also perform canonical quantization of the theory as required for the definition of the (physical) Hilbert space.
The key idea is to cure the above problem of non-invertibility of the kinetic term by removing the redundancy in the space of gauge field configurations due to gauge invariance. Let

$$
\begin{equation*}
\mathcal{A}:=\text { space of all gauge fields } A_{\mu}(x) \tag{9.30}
\end{equation*}
$$

and denote by

$$
\begin{equation*}
A_{\mu}^{h}(x):=h A_{\mu} h^{-1}+\frac{i}{g}\left(\partial_{\mu} h\right) h^{-1} \tag{9.31}
\end{equation*}
$$

a gauge transformation on $A_{\mu}^{h}$ with $h \in H$. Then fields related by a gauge transformation are physically equivalent, $A_{\mu}(x) \sim A_{\mu}^{h}(x)$. This is because $A_{\mu}$ satisfies the equation of motion if and only if so does $A_{\mu}^{h}$.
Given $A_{\mu}(x)$, the gauge equivalent field configurations are said to lie in the same orbit under gauge
transformations $h \in H$, denoted

$$
\begin{equation*}
\left\{A_{\mu}^{h}(x) \forall h \in H\right\} \tag{9.32}
\end{equation*}
$$

The physically inequivalent degrees of freedom are then counted by the quotient space

$$
\begin{equation*}
\mathcal{A} / H:=\left\{A_{\mu} \sim A_{\mu}^{h}: A_{\mu} \in \mathcal{A}, h \in H\right\}, \tag{9.33}
\end{equation*}
$$

i.e. they are given by all $A_{\mu}$-field configurations modulo gauge transformations. Our aim is now to define the path integral not on $\mathcal{A}$, but on the physical space $\mathcal{A} / H$ as

$$
\begin{equation*}
\int_{\mathcal{A} / H} \mathrm{~d} \mu[A] e^{i S[A]} . \tag{9.34}
\end{equation*}
$$

To this end we must construct a suitable measure $\mu[A]$ on $\mathcal{A} / H$ such that

$$
\begin{equation*}
\int_{\mathcal{A}} \mathcal{D} A=\int_{H} \mathrm{~d} \mu[h] \int_{\mathcal{A} / H} \mathrm{~d} \mu[A], \tag{9.35}
\end{equation*}
$$

with $\mathrm{d} \mu[h]$ a measure on $H$. If we have found such a measure, we can factor out the integration over the gauge orbit and only pick one representative per orbit as shown below.


To achieve (9.35) one introduces a gauge fixing condition

$$
\begin{equation*}
F(A) \stackrel{!}{=} 0 \tag{9.36}
\end{equation*}
$$

for some function $F$ which maps the $\operatorname{Lie}(H)$-valued connection $\mathcal{A}$ to some other $\operatorname{Lie}(H)$-valued field. For example we could take $F(A):=\partial_{\mu} A^{\mu}$ or generalisations thereof to be discussed later. Ideally, we would like that given $A \in \mathcal{A}$

$$
\begin{equation*}
F\left(A^{h}\right)=0 \quad \text { for some unique } \quad h \in H \tag{9.37}
\end{equation*}
$$

such as to fix one representative per gauge orbit. ${ }^{2}$ Then on the lefthand side of (9.35) we can insert $\mathbb{1}$ given by

$$
\begin{equation*}
\mathbb{1}=\int_{H} \mathrm{~d} \mu[h] \delta\left[F\left(A^{h}\right)\right] M\left(A^{h}\right) \tag{9.38}
\end{equation*}
$$

[^46]where the Jacobian factor
\[

$$
\begin{equation*}
M\left(A^{h}\right)=\operatorname{det} \frac{\partial F\left(A^{h}\right)}{\partial h} \tag{9.39}
\end{equation*}
$$

\]

is needed in analogy with the formula

$$
\begin{equation*}
\mathbb{1}=\int_{\mathbb{R}} \mathrm{d} x \delta[f(x)]\left|\frac{\partial f(x)}{\partial x}\right| \tag{9.40}
\end{equation*}
$$

and will be evaluated momentarily. Then

$$
\begin{align*}
\int_{\mathcal{A}} \mathcal{D A} & =\int_{\mathcal{F}} \mathcal{D} A \int_{H} \mathrm{~d} \mu[h] \delta\left[F\left(A^{h}\right)\right] M\left(A^{h}\right) \\
& =\int_{H} \mathrm{~d} \mu[h] \int_{\mathcal{A}} \mathcal{D} A \delta\left[F\left(A^{h}\right)\right] M\left(A^{h}\right) . \tag{9.41}
\end{align*}
$$

For all that follows we make the non-trivial assumption that we can find a $\mathcal{D A}$ that satisfies

$$
\begin{equation*}
\mathcal{D A}=\mathcal{D} A^{h} . \tag{9.42}
\end{equation*}
$$

If no such measure exists, the gauge symmetry is said to be anomalous and no consistent quantization procedure can be found. We will discuss such anomalies later in this course. Given (9.42) we have

$$
\begin{equation*}
\int_{\mathcal{A}} \mathcal{D} A=\int_{H} \mathrm{~d} \mu[h] \int_{\mathcal{A}} \mathcal{D} A^{h} \delta\left[F\left(A^{h}\right)\right] M\left(A^{h}\right) . \tag{9.43}
\end{equation*}
$$

Relabelling $A^{h} \rightarrow A$ yields

$$
\begin{equation*}
\int_{\mathcal{A}} \mathcal{D} A=\int_{H} \mathrm{~d} \mu[h] \underbrace{\int_{\mathcal{A}} \mathcal{D} A \delta[F(A)] M(A)}_{\substack{\equiv \int_{\mathcal{A} / H} \mathrm{~d} \mu[A]}}, \tag{9.44}
\end{equation*}
$$

which defines our desired measure on $\mathcal{A} / H$. This procedure can be repeated with any gauge invariant integrand. In particular, since $S[A]=S\left[A^{h}\right]$ by classical gauge invariance, this results in

$$
\begin{equation*}
\int_{\mathcal{A}} \mathcal{D} A e^{i S[A]}=\int_{H} \mathrm{~d} \mu[h] \int_{\mathcal{A}} \mathcal{D} A \delta[F(A)] M(A) e^{i S[A]} . \tag{9.45}
\end{equation*}
$$

The factor $\int_{H} \mathrm{~d} \mu[h]$ is just an overall factor $\operatorname{Vol}(H)$ that will cancel against the same factor in the denominator in the expression of correlation functions and can thus be ignored. Thus we define the partition function

$$
\begin{equation*}
Z=\int_{\mathcal{A}} \mathcal{D} A \delta[F(A)] M(A) e^{i S[A]} \equiv \int_{\mathcal{A} / H} \mathrm{~d} \mu[A] e^{i S[A]} \tag{9.46}
\end{equation*}
$$

It remains to compute $M(A)$ : By assumption, there exists one $h=h_{0}$ such that $F\left[A^{h_{0}}\right]=0$. Note that the measure $\mathrm{d} \mu[h]$ on $H$ satisfies

$$
\begin{equation*}
\int_{H} \mathrm{~d} \mu[h] f(h)=\int_{H} \mathrm{~d} \mu[h] f(g h) \quad \forall g \in H \tag{9.47}
\end{equation*}
$$

for all functions $f$ on $H$. Applying this to (9.38) with $g=h_{0}$ gives

$$
\begin{equation*}
\mathbb{1}=\int_{H} \mathrm{~d} \mu[h] \delta\left[F\left(\tilde{A}^{h}\right)\right] M\left(\tilde{A}^{h}\right) \quad \text { with } \quad \tilde{A}=A^{h_{0}} . \tag{9.48}
\end{equation*}
$$

This receives contributions only near $h=i d$. Thus we can express the integral as one over infinitesimal group elements by expanding $h=e^{i \alpha}$ with $\alpha \in \operatorname{Lie}(H)$,

$$
\begin{equation*}
\mathbb{1}=\int_{H} \mathrm{~d} \mu[h] \delta\left[F\left(\tilde{A}^{h}\right)\right] M\left(\tilde{A}^{h}\right)=\int_{\operatorname{Lie}(H)} \mathrm{d} \alpha \delta\left[F\left(\tilde{A}^{\alpha}\right)\right] M\left(\tilde{A}^{\alpha}\right) \tag{9.49}
\end{equation*}
$$

for a suitably normalised measure $\mathrm{d} \alpha$. Here

$$
\begin{equation*}
M\left(\tilde{A}^{\alpha}\right)=\operatorname{det} \frac{\partial F\left(\tilde{A}^{\alpha}\right)}{\partial \alpha} \tag{9.50}
\end{equation*}
$$

with

$$
\begin{equation*}
\frac{\partial F\left(\tilde{A}^{\alpha}\right)}{\partial \alpha}=\frac{\partial F\left(\tilde{A}^{\alpha}\right)}{\partial \tilde{A}^{\alpha}} \frac{\partial}{\partial \alpha} \tilde{A}^{\alpha} . \tag{9.51}
\end{equation*}
$$

Given the form $\tilde{A}_{\mu}^{\alpha}=\tilde{A}_{\mu}-D_{\mu} \alpha$ of an infinitesimal gauge transformation, (9.51) results in

$$
\begin{equation*}
\frac{\partial F\left(\tilde{A}^{\alpha}\right)}{\partial \tilde{A}_{\mu}^{\alpha}}\left(-D_{\mu}\right) . \tag{9.52}
\end{equation*}
$$

Finally we can replace $\tilde{A}^{\alpha} \rightarrow A$ as we did before because $\mathcal{D A}$ and $S[A]$ are gauge invariant. Thus we end up with

$$
\begin{equation*}
Z \simeq \int_{\mathcal{A}} \mathcal{D} A \delta[F(A)] \operatorname{det} \Delta_{\mathrm{FP}} e^{i S[A]} \tag{9.53}
\end{equation*}
$$

with the Faddeev-Popov matrix

$$
\begin{equation*}
\Delta_{\mathrm{FP}}=-\frac{\partial F(A)}{\partial A_{\mu}} D_{\mu} . \tag{9.54}
\end{equation*}
$$

With the help of this gauge fixed expression of the path integral one can now compute expectation values of any gauge invariant operator $O(A)$ as

$$
\begin{equation*}
\langle\Omega| T O(A)|\Omega\rangle=\lim _{T \rightarrow \infty(1-i \epsilon)} \frac{\int_{\mathcal{A}} \mathcal{D} A \delta[F(A)] \operatorname{det} \Delta_{\mathrm{FP}} O(A) e^{i} \int_{-T}^{T} \mathrm{~d}^{4} x \mathcal{L}}{\int_{\mathcal{A}} \mathcal{D} A \delta[F(A)] \operatorname{det} \Delta_{\mathrm{FP}} e^{i \int_{-T}^{T} \mathrm{~d}^{4} x \mathcal{L}}} . \tag{9.55}
\end{equation*}
$$

Note that gauge invariance of $O(A)$ is crucial as otherwise the replacement $A^{h} \rightarrow A$ at intermediate steps is wrong.

### 9.3 Faddeev-Popov ghosts

For all further computations a couple of further technical steps are required which help us rewrite

$$
\begin{equation*}
\int \mathcal{D} A \delta[F(A)] \operatorname{det} \Delta_{\mathrm{FP}} e^{i S[A]}: \tag{9.56}
\end{equation*}
$$

1. First, we can rewrite the delta-functional implementing the gauge-fixing constraint as

$$
\begin{equation*}
\delta[F(A)]=\int \mathcal{D} B^{a}(x) e^{i \int \mathrm{~d}^{4} x B^{a}(x) F^{a}(A)} \tag{9.57}
\end{equation*}
$$

where $F(A)$ is a $\operatorname{Lie}(H)$-valued field, i.e. $F(A)=F^{a}(A) T^{a}$ and thus likewise $B(x)=B^{a}(x) T^{a}$ is a $\operatorname{Lie}(H)$-valued auxiliary field called Nakanishi-Lautrup field with suitably normalised measure $\mathcal{D} B$.
2. To rewrite $\operatorname{det} \Delta_{\mathrm{FP}}$ recall the general formula

$$
\begin{equation*}
\int \mathrm{d}^{n} \theta^{*} \mathrm{~d} \theta e^{\theta_{i}^{*} M_{i j} \theta_{j}}=\operatorname{det} M \tag{9.58}
\end{equation*}
$$

for $M_{i j} \in \mathbb{C}^{N, N}, \theta_{i}^{*}, \theta_{i}$ complex Grassmann variables. The Faddeev-Popov matrix is an operator acting on $\operatorname{Lie}(H)$-valued fields. So we replace

$$
\begin{equation*}
\theta_{i}^{*} \rightarrow \bar{c}(x)=\bar{c}^{a}(x) T^{a}, \quad \theta_{i} \rightarrow c(x)=c^{a}(x) T^{a}, \tag{9.59}
\end{equation*}
$$

with $\bar{c}(x), c(x)$ Grassmann-valued, $\operatorname{Lie}(H)$-valued fields, to write

$$
\begin{equation*}
\operatorname{det} \Delta_{\mathrm{FP}}=\int \mathcal{D} \bar{c} \mathcal{D} c e^{i \int \mathrm{~d}^{4} x \bar{c}^{a}(x)\left(\Delta_{\mathrm{FP}}(x)\right)^{a}} \tag{9.60}
\end{equation*}
$$

with suitably normalised measures $\mathcal{D} \bar{c}, \mathcal{D} c$.
Thus altogether

$$
\begin{equation*}
\int \mathcal{D} A \delta[F(A)] \operatorname{det} \Delta_{\mathrm{FP}} e^{i S[A]}=\int \mathcal{D} A \int \mathcal{D} B \int \mathcal{D} \bar{c} \int \mathcal{D} c e^{i S[A, B, \bar{c}, c]}, \tag{9.61}
\end{equation*}
$$

where

$$
\begin{equation*}
S[A, B, \bar{c}, c]=\int \mathrm{d}^{4} x\left(-\frac{1}{4}\left(F_{\mu \nu}^{a}\right)^{2}+B^{a}(x) F^{a}(A)+\bar{c}^{a}(x)\left(\Delta_{\mathrm{FP}} c(x)\right)^{a}\right) \tag{9.62}
\end{equation*}
$$

The Grassmann-valued fields $c^{a}(x)$ and $\bar{c}^{a}(x)$ are called Faddeev-Popov ghosts and anti-ghosts respectively. They transform as scalar fields under $\operatorname{SO}(1,3)$, but have fermionic statistics due to their Grassmannian nature. Thus they violate the spin-statistics theorem. The underlying reason is that they violate unitarity since their Fock space does not have a positive definite norm, as will be discussed in detail later. This explains the name 'ghost', which generally describes a state with non-positive norm. We shall think of the ghost fields as unphysical, 'negative' degrees of freedom which, as we will also see explicitly later, cancel unphysical polarisations of the Yang-Mills field.

Instead of requiring that $F(A)=0$ we could equally well require that

$$
\begin{equation*}
F(A(x))=f(x) \tag{9.63}
\end{equation*}
$$

### 9.3. FADDEEV-POPOV GHOSTS

for some $\operatorname{Lie}(H)$-valued field $f(x)$. In fact, one can even average over all possible $f(x)$ with an arbitrary measure - because all results are independent of how precisely we fix the gauge. A convenient way to do this is to replace

$$
\begin{equation*}
\delta[F(A)] \rightarrow \int \mathcal{D} f \delta[F(A)-f] e^{-\frac{i}{2 \xi} \int \mathrm{~d}^{4} x f^{a}(x) f^{a}(x)} \tag{9.64}
\end{equation*}
$$

for arbitrary $\xi$. Then (9.64) can alternatively be written as

$$
\begin{equation*}
e^{-\frac{i}{2 \xi} \int \mathrm{~d}^{4} x F^{a}(A(x)) F^{a}(A(x))} \tag{9.65}
\end{equation*}
$$

or, by introducing as before an auxiliary field $B(x)$,

$$
\begin{align*}
\int \mathcal{D} B \int & \mathcal{D} f e^{-\frac{i}{2 \xi} \int \mathrm{~d}^{4} x f^{a}(x) f^{a}(x)+i \int \mathrm{~d}^{4} x B^{a}(x)\left(F^{a}(x)-f^{a}(x)\right)} \\
& =\int \mathcal{D} B \int \mathcal{D} f e^{-\frac{i}{2 \xi} \int \mathrm{~d}^{4} x\left(f^{a}(x)+\xi B^{a}(x)\right)^{2}+\frac{i}{2} \xi \int \mathrm{~d}^{4} x B^{a}(x) B^{a}(x)+i \int \mathrm{~d}^{4} x B^{a}(x) F^{a}(x)}  \tag{9.66}\\
& \simeq \int \mathcal{D} B e^{i \int d^{4} x\left(\frac{\xi}{2} B^{a}(x) B^{a}(x)+B^{a}(x) F^{a}(x)\right)},
\end{align*}
$$

where in the last step we are omitting an overall numerical factor from the Gaussian integral that will cancel in all correlation functions. Furthermore we abbreviate $F(A(x))^{a} \equiv F^{a}(x)$ (not to be confused with the Yang-Mills field strength $F_{\mu \nu}^{a}(x)$ ). Thus there are two alternative final results for the gauge fixed path integral:

- Using (9.65) we find

$$
\begin{equation*}
Z=\int \mathcal{D} A \int \mathcal{D} \bar{c} \int \mathcal{D} c e^{i \int \mathrm{~d}^{4} x\left(-\frac{1}{4}\left(F_{\mu \nu}^{a}\right)^{2}-\frac{1}{2 \xi} F^{a} F^{a}+\bar{c}^{a}\left(\Delta_{\mathrm{FP}}\right)^{a}\right]} \tag{9.67}
\end{equation*}
$$

- Using (9.66) yields

$$
\begin{equation*}
Z=\int \mathcal{D} A \int \mathcal{D} B \int \mathcal{D} \bar{c} \int \mathcal{D} c e^{i S[A, B, \bar{c}, c]} \tag{9.68}
\end{equation*}
$$

with

$$
\begin{equation*}
S[A, B, \bar{c}, c]=\int \mathrm{d}^{4} x\left(-\frac{1}{4}\left(F_{\mu \nu}^{a}\right)^{2}+B^{a} F^{a}+\frac{\xi}{2} B^{a} B^{a}+\bar{c}^{a}\left(\Delta_{\mathrm{FP}} c\right)^{a}\right) . \tag{9.69}
\end{equation*}
$$

Note that in (9.68) $B^{a}(x)$ has no kinetic term (hence its name auxiliary field). Its equation of motion is

$$
\begin{equation*}
B^{a}(x)=-\frac{1}{\xi} F^{a}(A(x)) \tag{9.70}
\end{equation*}
$$

Integrating out $B^{a}(x)$, i.e. replacing $B^{a}$ by its equation of motion and omitting $\int \mathcal{D} B$, gives (9.67) from (9.68).

In the sequel we will often work with the gauge-fixing condition

$$
\begin{equation*}
F(A)=\partial_{\mu} A^{\mu} \tag{9.71}
\end{equation*}
$$

and

$$
\begin{equation*}
\Delta_{\mathrm{FP}}=-\frac{\partial F}{\partial A_{\mu}} D_{\mu}=-\partial^{\mu} D_{\mu} \tag{9.72}
\end{equation*}
$$

Then e.g. (9.67) becomes

$$
\begin{equation*}
Z=\int \mathcal{D} A \int \mathcal{D} \bar{c} \int \mathcal{D} c e^{i S[A, c, \bar{c}]} \tag{9.73}
\end{equation*}
$$

with

$$
\begin{equation*}
S[A, \bar{c}, c]=\int \mathrm{d}^{4} x\left[-\frac{1}{4}\left(F_{\mu \nu}^{a}\right)^{2}-\frac{1}{2 \xi}\left(\partial_{\mu} A^{a \mu}\right)^{2}+\bar{c}^{a}\left(-\partial^{\mu} D_{\mu} c\right)^{a}\right] \tag{9.74}
\end{equation*}
$$

The value of the so-called gauge-fixing parameter $\xi$ is arbitrary, and we will see that all physical amplitudes are independent of this gauge choice.

### 9.4 Canonical quantisation and asymptotic Fock space

While the path-integral allows for the computation of $\langle\Omega| T O(A)|\Omega\rangle$ for any gauge-invariant operator $O(A)$ - we will provide the Feynman rules soon - canonical quantisation is required to establish the Hilbert space of states. As stressed in section (9.2), naive canonical quantization of the gauge invariant Yang-Mills Lagrangian

$$
\begin{equation*}
\mathcal{L}_{\mathrm{YM}}=-\frac{1}{4}\left(F_{\mu \nu}^{a}\right)^{2} \tag{9.75}
\end{equation*}
$$

faces the problem that $A_{0}^{a}$ has no canonically conjugate momentum. What saves the day is that the Faddeev-Popov treatment of the path-integral has provided us with a gauge-fixed classical Lagrangian

$$
\begin{equation*}
\mathcal{L}(A, B, \bar{c}, c)=-\frac{1}{4}\left(F_{\mu \nu}^{a}\right)^{2}-\partial^{\mu} B^{a} A_{\mu}^{a}+\frac{\xi}{2} B^{a} B^{a}+\partial^{\mu} \bar{c}^{a} D_{\mu} c^{a} \tag{9.76}
\end{equation*}
$$

Indeed our approach will be to view (9.76) as the proper Lagrangian that should serve as our starting point also for canonical quantization.
Starting from (9.76) gives the canonically conjugate fields

$$
\begin{align*}
& \Pi_{A}^{a \mu}=\frac{\partial \mathcal{L}}{\partial \dot{A}_{\mu}^{a}}=\left(F^{a}\right)^{\mu 0}, \quad \Pi_{c}^{a}=\frac{\mathcal{L}}{\partial \dot{c}^{a}}=-\dot{\bar{c}}^{a}  \tag{9.77}\\
& \Pi_{B}^{a}=\frac{\partial \mathcal{L}}{\partial \dot{B}^{a}}=-A_{0}^{a}, \quad \Pi_{\bar{c}}^{a}=\frac{\partial \mathcal{L}}{\partial \dot{\bar{c}}^{a}}=\dot{c}^{a}-g f^{a b c} A_{0}^{b} c^{c} \tag{9.78}
\end{align*}
$$

At first sight it looks as if the original problem that $\Pi_{A}^{a 0} \equiv 0$ in the non-gauge-fixed $\mathcal{L}_{\mathrm{YM}}$ still persists, after all $F^{a 00}$ still vanishes and it is only for $A_{i}^{a}$ that we find non-zero commutation relations

$$
\begin{equation*}
\left[A_{j}^{a}(t, \vec{x}), \dot{A}_{k}^{b}(t, \vec{y})\right]=i \delta^{a b} \delta_{k j} \delta^{(3)}(\vec{x}-\vec{y}) . \tag{9.79}
\end{equation*}
$$

These follow from

$$
\begin{equation*}
\left[A_{j}^{a}(t, \vec{x}), \Pi_{k}^{b}(t, \vec{y})\right]=\left[A_{j}^{a}(t, \vec{x}), \partial_{k} A_{0}^{b}(t, \vec{y})-\dot{A}_{k}^{b}(t, \vec{y})\right]=i \eta_{j k} \delta^{a b} \delta^{(3)}(\vec{x}-\vec{y}) \tag{9.80}
\end{equation*}
$$

together with $\left[A_{j}^{a}(t, \vec{x}), A_{k}^{a}(t, \vec{y})\right]=0$.

On the other hand (9.78) implies that

$$
\begin{equation*}
\left[B^{a}(t, \vec{x}),-A_{0}^{a}(t, \vec{y})\right]=i \delta^{a b} \delta^{(3)}(\vec{x}-\vec{y}) \tag{9.81}
\end{equation*}
$$

Thus we can view this as the missing commutation relation for $A_{0}^{a}$ by effectively treating $B^{a}$ as the momentum canonically conjugate to $A_{0}^{a}$ :

$$
\begin{equation*}
\left[A_{0}^{a}(t, \vec{x}), B^{b}(t, \vec{y})\right]=i \delta^{a b} \delta^{(3)}(\vec{x}-\vec{y}) . \tag{9.82}
\end{equation*}
$$

By the Grassmanian nature of ghosts we furthermore impose the anti-commutation relations

$$
\begin{equation*}
\left\{c^{a}(t, \vec{x}), \Pi_{c}^{b}(t, \vec{y})\right\}=\left\{\bar{c}^{a}(t, \vec{x}), \Pi_{\bar{c}}^{b}(t, \vec{y})\right\}=i \delta^{a b} \delta^{(3)}(\vec{x}-\vec{y}) . \tag{9.83}
\end{equation*}
$$

Note that with

$$
\begin{equation*}
c^{a \dagger}=c^{a}, \quad\left(\bar{c}^{a}\right)^{\dagger}=-\bar{c}^{a} \tag{9.84}
\end{equation*}
$$

we have $\mathcal{L}=\mathcal{L}^{\dagger}$ as required for unitarity of the full $S$-matrix.
In trying to determine the physical Hilbert space we face the problem that the theory is intrinsically interacting and thus the Heisenberg operators have no free-field mode expansion. The full spectrum therefore contains complicated bound states of gauge bosons, and, if present, of charged matter states. In fact, deriving the Hilbert space of Yang-Mills theory from first principles is one of the biggest open questions of Quantum Field Theory.
However, if we content ourselves with computing matrix elements for scattering experiments, it suffices to know the asymptotic in- and out-states. According to the general LSZ logic developed in section 2.3 , these correspond to the Fock space of excitations created from the vacuum by the asymptotic in- and out-fields. The asymptotic fields in turn are given by fully renormalized, but free fields whose dynamics is determined by the free Lagrangian (i.e. setting $g=0$ )

$$
\begin{equation*}
\tilde{\mathcal{L}}_{0}=-\frac{1}{4}\left(\partial_{\mu} \tilde{A}_{v}^{a}-\partial_{v} \tilde{A}_{\mu}^{a}\right)^{2}-\partial^{\mu} \tilde{B}^{a} \tilde{A}_{\mu}^{a}+\frac{\tilde{\xi}}{2} \tilde{B}^{a} \tilde{B}^{a}+\partial^{\mu} \tilde{c}^{a} \partial_{\mu} \tilde{c}^{a}, \tag{9.85}
\end{equation*}
$$

where the tilde stresses the nature as asymptotic fields. In the so-called Feynman gauge corresponding to $\tilde{\xi}=1$ all fields enjoy a free mode expansion (see the tutorials for a proof). In particular we can expand the asymptotic gauge field as

$$
\begin{equation*}
\tilde{A}_{\mu}^{a}(x)=\int \frac{\mathrm{d}^{3} k}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{k}}} \sum_{\lambda=0}^{3} \epsilon_{\mu}^{a}(\vec{k}, \lambda)\left[a_{\lambda}^{a}(\vec{k}) e^{-i k \cdot x}+a_{\lambda}^{a \dagger}(\vec{k}) e^{i k \cdot x}\right], \tag{9.86}
\end{equation*}
$$

where the polarisation vectors satisfy

$$
\begin{equation*}
\epsilon(\vec{k}, \lambda) \cdot \epsilon\left(\vec{k}, \lambda^{\prime}\right)=\eta_{\lambda \lambda^{\prime}} . \tag{9.87}
\end{equation*}
$$

Furthermore one can integrate out $\tilde{B}^{a}$ by replacing it by its equation of motion $\tilde{B}^{a}=-\partial^{\mu} \tilde{A}_{\mu}^{a}$ (for $\tilde{\xi}=1$ ) and combine (9.79) and (9.82) into the commutation relations

$$
\begin{equation*}
\left[A_{\mu}^{a}(t, \vec{x}), \dot{A}_{\nu}^{b}(t, \vec{y})\right]=-i \delta^{a b} \eta_{\mu \nu} \delta^{(3)}(\vec{x}-\vec{y}), \tag{9.88}
\end{equation*}
$$

i.e.

$$
\begin{equation*}
\left[a_{\lambda}^{a \dagger}(\vec{k}), a_{\lambda^{\prime}}^{b}\left(\vec{k}^{\prime}\right)\right]=-\delta^{a b} \eta_{\lambda \lambda^{\prime}}(2 \pi)^{3} \delta^{(3)}\left(\vec{k}-\vec{k}^{\prime}\right) \tag{9.89}
\end{equation*}
$$

As in QED the resulting asymptotic Fock spaces contains negative norm states for time-like and and zero-norm states for longitudinal polarisations. Furthermore

$$
\begin{equation*}
\tilde{c}^{a}(x)=\int \frac{\mathrm{d}^{3} p}{(2 \pi)^{3}} \frac{1}{\sqrt{2 E_{p}}}\left(c^{a}(\vec{p}) e^{-i p \cdot x}+c^{a \dagger}(\vec{p}) e^{i p \cdot x}\right) \tag{9.90}
\end{equation*}
$$

(with a similar expansion for $\bar{c}$ ), and the anti-commutation relations for the asymptotic ghost modes read

$$
\begin{equation*}
\left\{c^{a}(\vec{p}), \bar{c}^{b \dagger}\left(\vec{p}^{\prime}\right)\right\} \text { and }\left\{\bar{c}^{a}(\vec{p}), c^{b \dagger}\left(\vec{p}^{\prime}\right)\right\} \propto-\delta^{a b} \delta^{(3)}\left(\vec{p}-\vec{p}^{\prime}\right) \tag{9.91}
\end{equation*}
$$

This yields zero-norm excitations in the ghost sector (see again the tutorials for more details).
In view of the appearance of negative and zero-norm states in the asymptotic Fock space - not to mention the problem that the ghost fields are non-physical degrees of freedom in the first place - the question arises: What is the physical Hilbert space?

### 9.5 BRST symmetry and the physical Hilbert space

The problem of determining the physical Hilbert space $\mathcal{H}_{\text {phys }}$ comes in two parts:

1. Give a criterion for $\mathcal{H}_{\text {phys }}$ which guarantees a positive-definite norm on $\mathcal{H}_{\text {phys }}$.
2. Show that time-evolution does not lead out of $\mathcal{H}_{\text {phys }}$, i.e. show that the $S$-matrix acts as a unitary operator on $\mathcal{H}_{\text {phys }}$.

In principle we are free to choose any criterion for $\mathcal{H}_{\text {phys }}$ as long as it leads to a positive-definite norm - after all this is part of the definition of the theory. However, in almost all cases the S-matrix will not respect this criterion. On the other hand, if the criterion for $\mathcal{H}_{\text {phys }}$ is related to a symmetry of the full interacting theory, then invariance of the physical Hilbert space under time-evolution follows because the $S$-matrix respects that symmetry by assumption.
As a preparation for our definition of $\mathcal{H}_{\text {phys }}$ we must therefore first determine the symmetries of the the full interacting, gauge-fixed Lagrangian

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu}{ }^{a} F^{\mu \nu a}+\frac{\xi}{2} B^{a} B^{a}+B^{a} \partial^{\mu} A_{\mu}^{a}+\bar{c}^{a}\left(-\partial^{\mu} D_{\mu} c\right)^{a} \tag{9.92}
\end{equation*}
$$

The first term is gauge invariant, but the gauge fixing term $B^{a} \partial^{\mu} A_{\mu}^{a}$ and the ghost Lagrangian $\sim \partial^{\mu} \bar{c}^{D} D_{\mu} c$ break gauge invariance as a result of the partial derivatives. This is of course as it must be - after all our desire to fix the gauge was what lead to the above Lagrangian.
However (9.92) possesses a residual global fermionic symmetry - the BRST symmetry ${ }^{3}$ - acting as $^{2}$ follows: First define a Grassmann odd operator $S$ such that

$$
\begin{equation*}
S A_{\mu}=-D_{\mu} c=-\left(\partial_{\mu} c+i g\left[A_{\mu}, c\right]\right), \quad S c=\frac{i}{2} g\{c, c\}, \quad S \bar{c}=-B, \quad S B=0 \tag{9.93}
\end{equation*}
$$

[^47]Note that

- $S A_{\mu}$ is just an infinitesimal gauge transformation $\delta A$ of $A_{\mu}$, but with $\delta A_{\mu}=-D_{\mu} \alpha$ replaced by $-D_{\mu} c$;
- the assertion that $S$ is Grassmann odd means that, by definition, $S$ obeys a graded Leibniz rule, e.g.

$$
\begin{equation*}
S(c A)=(S c) A-c(S A), \tag{9.94}
\end{equation*}
$$

since the $c$ is Grassmann odd.
The operator $S$ has two decisive properties:

1. $S$ is nilpotent,

$$
\begin{equation*}
S^{2}=0, \tag{9.95}
\end{equation*}
$$

i.e. $S \cdot S \cdot \Phi=0$ with $\Phi$ any of the fields $A, c, \bar{c}, B$.
2. The Lagrangian (9.92) can be written as

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4} F_{\mu \nu}{ }^{a} F^{\mu \nu a}-S \psi(x) \tag{9.96}
\end{equation*}
$$

for $\psi(x)=\bar{c}^{a} \partial^{\mu} A_{\mu}^{a}+\frac{\xi}{2} \bar{c}^{a} B^{a}$.
Note that properties (9.95) and (9.96) immediately imply that the action (9.92) is invariant under application of $S$, i.e.

$$
\begin{equation*}
S \cdot \mathcal{L}=0 \tag{9.97}
\end{equation*}
$$

because $S \cdot\left(F_{\mu \nu}{ }^{a} F^{\mu \nu a}\right)=0$ since $S$ acting on $A_{\mu}$ is just a gauge transformation with gauge parameter $c(x)$ and by (9.95) $S^{2} \cdot \psi=0$.
The proof of (9.95) is worked out on the exercise sheets and requires the Jacobi identity. Property (9.96) follows simply by noting that

$$
\begin{align*}
S \cdot\left(\bar{c}^{a} \partial^{\mu} A_{\mu}^{a}\right) & =(S \bar{c})^{a}\left(\partial^{\mu} A_{\mu}\right)^{a}-\bar{c}^{a} \partial^{\mu}\left(S A_{\mu}\right)^{a}=-B^{a}\left(\partial_{\mu} A^{\mu}\right)^{a}+\bar{c}^{a} \partial^{\mu}\left(D_{\mu} c\right)^{a},  \tag{9.98}\\
S \frac{\xi}{2} \bar{c}^{a} B^{a} & =\frac{\xi}{2}\left((S \bar{c})^{a} B^{a}-\bar{c}^{a} S B^{a}\right)=\frac{\xi}{2}\left(-B^{a} B^{a}-\bar{c}^{a} \cdot 0\right) . \tag{9.99}
\end{align*}
$$

Finally, the BRST symmetry transformation is defined as

$$
\begin{equation*}
\delta_{\epsilon} \Phi=\epsilon S \cdot \Phi \tag{9.100}
\end{equation*}
$$

with $\Phi$ any of the fields $A, c, \bar{c}, B$. Here $\epsilon$ is a global (i.e. not spacetime-dependent) Grassmann variable. Clearly $\delta_{\epsilon}^{2}=0$ because $S^{2}=0$ and

$$
\begin{equation*}
\delta_{\epsilon} \mathcal{L}=0 \tag{9.101}
\end{equation*}
$$

By Noether's theorem one can construct the associated classical Noether current $J_{\text {BRST }}^{\mu}$ with

$$
\begin{equation*}
\partial_{\mu} J_{\mathrm{BRST}}^{\mu}=0 \tag{9.102}
\end{equation*}
$$

and its charge

$$
\begin{equation*}
Q_{\mathrm{BRST}}=\int \mathrm{d}^{4} x J_{\mathrm{BRST}}^{0} \quad \text { with } \quad \dot{Q}_{\mathrm{BRST}}=0 \tag{9.103}
\end{equation*}
$$

The precise form of $J_{\text {BRST }}^{\mu}$ and $Q_{\text {BRST }}$ can easily be worked out, but is not essential for us here. If we perform canonical quantisation as sketched in the previous section, the classical Noether charge $Q_{\text {BRST }}$ is promoted to a quantum BRST operator $\hat{Q}_{\text {BRST }}$ which generates the BRST symmetry transformation in the sense that

$$
\begin{equation*}
\left[\epsilon \hat{Q}_{\mathrm{BRST}}, \hat{X}\right]=i \delta_{\epsilon} \hat{X} . \tag{9.104}
\end{equation*}
$$

That a Noether charge generates the symmetry underlying its conservation in the quantum theory is a general fact. (9.104) can also be checked explicitly via the canonical commutation relations once $\hat{Q}_{B R S T}$ is worked out.
As a result of the nilpotence of $S$, also $\hat{Q}_{\text {BRST }}$ is nilpotent,

$$
\begin{equation*}
\hat{Q}_{\mathrm{BRST}}^{2}=0 \tag{9.105}
\end{equation*}
$$

In order to see (9.104) we first consider $\hat{X}$ to be bosonic. Then

$$
\begin{equation*}
[\epsilon \hat{Q}, \hat{X}]=i \delta_{\epsilon} \hat{X}=i \epsilon S \hat{X} \tag{9.106}
\end{equation*}
$$

implies

$$
\begin{equation*}
[\hat{Q}, \hat{X}]=i S \hat{X} \tag{9.107}
\end{equation*}
$$

If $\hat{X}$ is fermionic, then instead

$$
\begin{equation*}
i \epsilon S \hat{X}=[\epsilon \hat{Q}, \hat{X}]=\epsilon \hat{Q} \hat{X}-\hat{X} \epsilon \hat{Q}=\epsilon(\hat{Q} \hat{X}+\hat{X} \hat{Q}) \tag{9.108}
\end{equation*}
$$

so $\{\hat{Q}, \hat{X}\}=i S \hat{X}$. Let $\hat{X}$ be bosonic. Then $S \hat{X}=-i[\hat{Q}, \hat{X}]$ is fermionic and

$$
\begin{equation*}
0=S^{2} \hat{X}=\{\hat{Q},[\hat{Q}, \hat{X}]\}=\hat{Q}(\hat{Q} \hat{X}-\hat{X} \hat{Q})+(\hat{Q} \hat{X}-\hat{X} \hat{Q}) \hat{Q}=\left[\hat{Q}^{2}, \hat{X}\right] . \tag{9.109}
\end{equation*}
$$

A similar conclusion holds for $\hat{X}$ fermionic, and thus

$$
\begin{equation*}
0=\left[\hat{Q}^{2}, \hat{X}\right] \text { for all fields } \tag{9.110}
\end{equation*}
$$

So either $\hat{Q}^{2} \sim \mathbb{1}$ or $\hat{Q}^{2}=0$. But $\hat{Q}^{2} \sim \mathbb{1}$ is not possible because symmetry of the theory under BRST transformations implies that also the vacuum must be invariant,

$$
\begin{equation*}
\hat{Q}|\Omega\rangle=0, \tag{9.111}
\end{equation*}
$$

in contradiction with $\hat{Q}^{2} \sim \mathbb{1}$.

Conservation of the classical Noether charge, $\dot{Q}_{\text {BRST }}=0$, translates into the relation

$$
\begin{equation*}
\left[\hat{H}, \hat{Q}_{\mathrm{BRST}}\right]=0 . \tag{9.112}
\end{equation*}
$$

Finally we note that

$$
\begin{equation*}
\hat{Q}_{\mathrm{BRST}}^{\dagger}=\hat{Q}_{\mathrm{BRST}}, \tag{9.113}
\end{equation*}
$$

as follows from the explicit expression for $Q_{\text {BRST }}$ together with

$$
\begin{equation*}
\hat{A}=\hat{A}^{\dagger}, \quad \hat{c}=\hat{c}^{\dagger}, \quad \hat{c}=-\hat{c}^{\dagger}, \quad \hat{B}=\hat{B}^{\dagger}, \tag{9.114}
\end{equation*}
$$

even though we will not derive this here.
Before discussing the relation of $\hat{Q}_{\text {BRST }}$ with the physical Hilbert space $\mathcal{H}_{\text {phys }}$, we embark on a little mathematical digression:
Let $\mathcal{H}$ be a vector space and $\hat{Q}: \mathcal{H} \mapsto \mathcal{H}$ a general nilpotent linear operator, i.e. $\hat{Q}^{2}=0$.

- An element $|\psi\rangle \in \mathcal{H}$ with $\hat{Q}|\psi\rangle=0$, i.e. $|\psi\rangle \in \operatorname{ker}(\hat{Q})$ is called $\hat{Q}$-closed.
- An element $|\psi\rangle \in \mathcal{H}$ such that $\exists \chi \in \mathcal{H}$ with $|\psi\rangle=\hat{Q}|\chi\rangle$, i.e. $|\psi\rangle \in \operatorname{Im} \hat{Q}$, is called $\hat{Q}$-exact.
- Since $\hat{Q}^{2}=0$, we have that

$$
\begin{equation*}
\operatorname{Im}(\hat{Q}) \subset \operatorname{Ker}(\hat{Q}) \tag{9.115}
\end{equation*}
$$

because $\hat{Q}|\psi\rangle=\hat{Q}^{2}|\chi\rangle=0$.

- If $\mathcal{H}$ has an inner product and if $\hat{Q}=\hat{Q}^{\dagger}$ with respect to it, then

$$
\begin{equation*}
\langle\psi \mid \psi\rangle=0 \quad \forall \psi \in \operatorname{Im}(\hat{Q}) \tag{9.116}
\end{equation*}
$$

because

$$
\begin{equation*}
\langle\psi \mid \psi\rangle=\langle\chi| \hat{Q}^{\dagger} \hat{Q}|\chi\rangle=\langle\chi| \hat{Q}^{2}|\chi\rangle=0 . \tag{9.117}
\end{equation*}
$$

Furthermore

$$
\begin{equation*}
\langle\psi \mid \phi\rangle=0 \quad \forall|\psi\rangle \in \operatorname{Im}(\hat{Q}),|\phi\rangle \in \operatorname{Ker}(\hat{Q}) . \tag{9.118}
\end{equation*}
$$

Thus $\hat{Q}$-exact states are null and orthogonal to $\hat{Q}$-closed states.

- As a result, $\hat{Q}$-closed states differing only by $\hat{Q}$-exact elements have the same norm. To see this let $\left|\psi_{1}\right\rangle,\left|\psi_{2}\right\rangle \in \operatorname{Ker}(\hat{Q})$ such that

$$
\begin{equation*}
\left|\psi_{1}\right\rangle-\left|\psi_{2}\right\rangle=\hat{Q}|\chi\rangle . \tag{9.119}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left\langle\psi_{1} \mid \psi_{1}\right\rangle=\left\langle\psi_{2} \mid \psi_{2}\right\rangle+\langle\chi| \hat{Q}^{\dagger} \hat{Q}|\chi\rangle+\left\langle\psi_{2} \mid \hat{Q} \chi\right\rangle+\langle\chi| \hat{Q}^{\dagger}\left|\psi_{2}\right\rangle=\left\langle\psi_{2} \mid \psi_{2}\right\rangle . \tag{9.120}
\end{equation*}
$$

- This motivates defining an equivalence relation on $\operatorname{Ker}(\hat{Q})$ by identifying elements in $\operatorname{Ker}(\hat{Q})$ which differ only by elements in $\operatorname{Im}(\hat{Q})$, i.e.

$$
\begin{equation*}
\forall \psi_{1}, \psi_{2} \in \operatorname{Ker}(\hat{Q}):\left|\psi_{1}\right\rangle \sim\left|\psi_{2}\right\rangle \quad \text { if } \quad\left|\psi_{1}\right\rangle-\left|\psi_{2}\right\rangle \in \operatorname{Im}(\hat{Q}) . \tag{9.121}
\end{equation*}
$$

One furthermore defines the $\hat{Q}$-cohomology as the quotient space

$$
\begin{equation*}
\operatorname{Cohom}(\hat{Q}):=\frac{\operatorname{Ker}(\hat{Q})}{\operatorname{Im}(\hat{Q})} \equiv \frac{\{\text { closed }|\psi\rangle\}}{\{\text { exact }|\psi\rangle\}} . \tag{9.122}
\end{equation*}
$$

- As an example of this structure consider $\mathcal{H}=\Omega^{\bullet}(M):=\{$ space of differential forms on manifold $M\}$ and $\hat{Q}:=d \wedge$, i.e. the exterior derivative. A closed $p$-form $\omega$ satisfies $d \omega=0$. An exact $p$-form $\chi$ satisfies $\chi=d \phi$ for some $(p-1)$-form, and the $p$-th cohomology group is

$$
\begin{equation*}
H^{p}(M):=\frac{\{\text { closed } p \text {-forms }\}}{\{\text { exact } p \text {-forms }\}} \tag{9.123}
\end{equation*}
$$

After this digression, we come back to our physical problem. In order to fully define quantum Yang-Mills-theory, we must define what we mean by the physical state space $\mathcal{H}_{\text {phy }}$, and then prove that this definition is stable under time-evolution and implies a positive-norm Hilbert space.
Our definition of $\mathcal{H}_{\text {phys }}$ is guided by the following two observations:

1. The space of states on which the time-evolution operator is independent of the specific choice of gauge-fixing condition is given by $\operatorname{Ker}\left(\hat{Q}_{\mathrm{BRST}}\right)$, i.e. by $|\psi\rangle$ such that $\hat{Q}_{\mathrm{BRST}}|\psi\rangle=0$. We will prove this below.
2. Within this space, the BRST-exact states $|\chi\rangle \in \operatorname{Im}\left(\hat{Q}_{\text {BRST }}\right)$, have zero overlap with all other states - as follows from previous discussion, in particular (9.118).

Now, invariance of time-evolution under changes of gauge fixing condition is obviously required for the theory to make sense, so we require for sure

$$
\begin{equation*}
\hat{Q}_{\text {BRST }}\left|\psi_{\text {phys }}\right\rangle \stackrel{!}{=} 0 \tag{9.124}
\end{equation*}
$$

for physical states. Furthermore, states with zero overlap with all other states can never be measured and are therefore unphysical. With this motivation we define

$$
\begin{equation*}
\mathcal{H}_{\text {phys }}=\operatorname{Cohom}\left(\hat{Q}_{\mathrm{BRST}}\right) \equiv \frac{\operatorname{Ker}\left(\hat{Q}_{\mathrm{BRST}}\right)}{\operatorname{Im}\left(\hat{Q}_{\mathrm{BRST}}\right)} . \tag{9.125}
\end{equation*}
$$

Before exploring $\mathcal{H}_{\text {phys }}$ further, we prove 1. from above: The time-evolution operator is

$$
\begin{equation*}
\hat{U}=T e^{i \int \hat{H} \mathrm{~d} t}=e^{i \int(\Pi \dot{\phi}-\mathcal{L}) \mathrm{d} t}, \tag{9.126}
\end{equation*}
$$

where $\Pi \dot{\phi}$ stands symbolically for all fields and their conjugates and $\mathcal{L}$ can be written as

$$
\begin{equation*}
\mathcal{L}=-\frac{1}{4}\left(F_{\mu \nu}^{a}\right)^{2}-S \cdot \psi, \tag{9.127}
\end{equation*}
$$

with $S$ generating the BRST transformations. Now, only $\psi$ depends on the explicit choice of gauge fixing condition $F(A)=0$. So a general transition amplitude $\langle\alpha| \hat{U}|\beta\rangle \simeq\langle\alpha|\left(1+i \int(\Pi \dot{\phi}-\mathcal{L})|\beta\rangle\right.$ (infinitesimally) changes under a change of the gauge fixing condition $F$ as

$$
\begin{equation*}
\delta_{F}\langle\alpha| \hat{U}|\beta\rangle=\langle\alpha|\left(i \delta_{F} S \psi\right)|\beta\rangle=\langle\alpha| i S \delta_{F} \psi|\beta\rangle . \tag{9.128}
\end{equation*}
$$

Since by definition iS $\hat{X}=\left[\hat{Q}_{\mathrm{BRST}}, \hat{X}\right]$, this is

$$
\begin{equation*}
=\langle\alpha|\left(\hat{Q} \delta_{F} \psi-\delta_{F} \psi \hat{Q}\right)|\beta\rangle . \tag{9.129}
\end{equation*}
$$

For this to vanish for all $\delta_{F} \psi$, we need

$$
\begin{equation*}
\hat{Q}_{\text {BRST }}|\beta\rangle=0=\langle\alpha| \hat{Q}_{\text {BRST }} . \tag{9.130}
\end{equation*}
$$

It remains to evaluate the physical state condition explicitly and to check the existence of a positivedefinitene norm on $\mathcal{H}_{\text {phys }}$, as well as to check unitarity of the $S$-matrix.

## Evaluating the physical state condition

We can only do this explicitly for the asymptotic in/out-states, for which we now how they are created by the action of the free asymptotic in/out-fields $\tilde{A}, \tilde{B}, \tilde{c}, \tilde{c}$. Since $\tilde{B} \sim \partial \tilde{A}$ by its equation of motion it suffices to consider 1-particle excitations of $\tilde{A}^{\mu}, \tilde{c}, \tilde{c}$. Let us define the states

$$
\begin{align*}
& \left|\tilde{A}_{\mu}^{a}(\vec{k})\right\rangle:=\sum_{\lambda=0}^{3} \epsilon_{\mu}^{a}(\vec{k}, \lambda) a_{\lambda}^{\dagger a}(\vec{k})|\Omega\rangle,  \tag{9.131}\\
& \left|\tilde{c}^{a}(\vec{k})\right\rangle:=c^{a \dagger}(\vec{k})|\Omega\rangle
\end{align*}
$$

and similarly for $\left|\tilde{c}^{a}(\vec{k})\right\rangle$, with $a^{\dagger}, c^{\dagger}, \bar{c}^{\dagger}$ the creation operators for the asymptotic fields $\tilde{A}, \tilde{c}, \tilde{\bar{c}}$, as introduced in section 9.4. Now, for these free asymptotic fields (for which essentially $g \equiv 0$ ), the BRST transformations take the simple form

$$
\begin{align*}
S \tilde{A}_{\mu} & =-\partial_{\mu} \tilde{c}, \\
S \tilde{c} & =0, \\
S \tilde{c} & =-\tilde{B}=-\partial \tilde{A},  \tag{9.132}\\
S \tilde{B} & =0,
\end{align*}
$$

in Feynman gauge $\tilde{\xi}=1$. Thus

$$
\begin{align*}
\hat{Q}_{\mathrm{BRST}}\left|\tilde{A}^{\mu a}(\vec{k})\right\rangle & =\alpha k^{\mu}\left|\tilde{c}^{a}(\vec{k})\right\rangle,  \tag{9.133}\\
\hat{Q}_{\mathrm{BRST}}\left|\tilde{c}^{a}(\vec{k})\right\rangle & =0,  \tag{9.134}\\
\hat{Q}_{\mathrm{BRST}}\left|\tilde{\bar{c}}^{a}(\vec{k})\right\rangle & =\beta k^{\mu}\left|\tilde{A}_{\mu}^{a}(\vec{k})\right\rangle \tag{9.135}
\end{align*}
$$

for some numerical factors $\alpha, \beta \neq 0$. First, we note that nil potency $\hat{Q}_{\mathrm{BRST}}^{2}=0$ implies that

$$
\begin{equation*}
0=\hat{Q}_{\mathrm{BRST}}^{2}\left|\bar{c}^{a}(k)\right\rangle=\alpha \beta\left(k^{\mu} k_{\mu}\right)\left|c^{a}(\vec{k})\right\rangle \tag{9.136}
\end{equation*}
$$

and thus

$$
\begin{equation*}
k^{2}=0 . \tag{9.137}
\end{equation*}
$$

## Properties of Cohom ( $\hat{Q}_{\text {BRST }}$ ):

- First evaluate the condition $\hat{Q}_{\text {BRST }}|\psi\rangle=0$ if $|\psi\rangle \in \mathcal{H}_{\text {phys }}$ : Applied to (9.135) this implies that a physical state $|\psi\rangle$ (with non-vanishing $k^{\mu}$ ) cannot have any $\bar{c}$-excitations. Furthermore, in view of (9.133), a state $|\xi\rangle:=\xi_{\mu}\left|A^{\mu}(\vec{k})\right\rangle$ with polarisation vector $\xi^{\mu}$ is physical only if

$$
\begin{equation*}
\xi^{\mu} k_{\mu}=0 . \tag{9.138}
\end{equation*}
$$

- Next analyse our identification $|\psi\rangle \sim 0$ if $|\psi\rangle=\hat{Q}_{\text {BRST }}|\chi\rangle$ : (9.133) implies $\left|\tilde{c}^{a}(\vec{k})\right\rangle$ is BRST trivial, while (9.135) implies that $\xi^{\mu}\left|A_{\mu}(k)\right\rangle$ is BRST trivial for $\xi^{\mu}=k^{\mu}$.

To summarise:
A physical state $|\psi\rangle \in \mathcal{H}_{\text {phys }}=\operatorname{Cohom}\left(\hat{Q}_{\mathrm{BRST}}\right)$ is of the form

$$
\begin{equation*}
|\psi\rangle=\xi^{\mu}\left|\tilde{A}_{\mu}(\vec{k})\right\rangle \tag{9.139}
\end{equation*}
$$

with $k^{2}=0$ for $\xi^{\mu} k_{\mu}=0$. We furthermore identify states with polarisation vector $\xi^{\mu} \sim \xi^{\mu}+k^{\mu}$, i.e. longitudinal null states are trivial.

## Unitarity of $S$-matrix:

Since $\mathcal{L}=\mathcal{L}^{\dagger}$, we know that the $S$-matrix is unitary on the entire Fock space with basis $|\gamma\rangle \in \mathcal{H}_{\text {Fock }}$, i.e.

$$
\begin{equation*}
\sum_{|\gamma\rangle}\langle\alpha| S^{\dagger}|\gamma\rangle\langle\gamma| S|\beta\rangle=\langle\alpha| \mathbb{1}|\beta\rangle . \tag{9.140}
\end{equation*}
$$

However, we need unitarity to hold on $\mathcal{H}_{\text {phys }}=\operatorname{Cohom}\left(\hat{Q}_{\text {BRST }}\right)$, i.e. ${ }^{4}$

$$
\begin{equation*}
\sum_{\left\langle\chi_{T}\right\rangle \in \mathcal{H}_{\text {phys }}}\left\langle\phi_{T}\right| S^{\dagger}\left|\chi_{T}\right\rangle\left\langle\chi_{T}\right| S\left|\psi_{T}\right\rangle \stackrel{!}{=}\left\langle\phi_{T}\right| \mathbb{1}\left|\psi_{T}\right\rangle \tag{9.141}
\end{equation*}
$$

for $\left|\phi_{T}\right\rangle,\left|\psi_{T}\right\rangle \in \mathcal{H}_{\text {phys }}$ transversely polarised states. Since $\left[\hat{H}, \hat{Q}_{\mathrm{BRST}}\right]=0$ we know that

$$
\begin{equation*}
\hat{Q}_{\mathrm{BRST}} S\left|\psi_{T}\right\rangle=S \hat{Q}_{\mathrm{BRST}}\left|\psi_{T}\right\rangle=0, \tag{9.142}
\end{equation*}
$$

so $S$ does not lead out of $\operatorname{Ker}(\hat{Q})$. But $S$ can and in fact does produce BRST exact states in $\operatorname{Im}\left(\hat{Q}_{\text {BRST }}\right)$ - we will see this explicitly in examples. Still, these are null, i.e. have zero overlap with all other states in $\operatorname{Ker}(\hat{Q})$. Therefore using (9.140)

$$
\begin{align*}
\left\langle\phi_{T}\right| \mathbb{1}\left|\psi_{T}\right\rangle & =\sum_{|\gamma\rangle}\left\langle\phi_{T}\right| S^{\dagger}|\gamma\rangle\langle\gamma| S\left|\psi_{T}\right\rangle \\
& =\sum_{\left\langle\chi_{T}\right\rangle \in \mathcal{H}_{\text {phys }}}\left\langle\phi_{T}\right| S^{\dagger}\left|\chi_{T}\right\rangle\left\langle\chi_{T}\right| S\left|\psi_{T}\right\rangle . \tag{9.143}
\end{align*}
$$

In the last line we first decomposed $\mathcal{H}_{\text {Fock }}$ into $\operatorname{Ker}\left(\hat{Q}_{\text {BRST }}\right)$ and its orthogonal complement and second used that only the projection of $|\gamma\rangle$ to the transversely polarized physical space Cohom $\left(\hat{Q}_{\text {BRST }}\right) \subset$ $\operatorname{Ker}\left(\hat{Q}_{\text {BRST }}\right)$ has non-trivial overlap with $S\left|\psi_{T}\right\rangle$ because the latter lies inside $\operatorname{Ker}\left(\hat{Q}_{\text {BRST }}\right)$, and states within $\operatorname{Im}\left(\hat{Q}_{B R S T}\right)$ are null and orthogonal on each other. Thus the $S$-matrix is unitary.

All of the above applies also to the special case of $U(1)$ gauge theory. There, the $\bar{c}-c$ system

[^48]does not couple to $A_{\mu}$ since $D_{\mu} c=\partial_{\mu} c$ as the ghosts transform in the adjoint representation, which for gauge group $U(1)$ is trivial. Also, the constraint
\[

$$
\begin{equation*}
Q_{\mathrm{BRST}}|\psi\rangle=0 \tag{9.144}
\end{equation*}
$$

\]

reduces to the Gupta-Bleuler condition

$$
\begin{equation*}
\partial A^{+}|\psi\rangle=0 \tag{9.145}
\end{equation*}
$$

In YM theory, imposing (9.140) on the full Hilbert space is not adequate since unlike the full operator $\hat{Q}_{\text {BRST }}$ it does not commute with the interacting Hamiltonian $\hat{H}$ and is thus not stable under timeevolution.


[^0]:    ${ }^{1}$ For corrections and improvement suggestions please send a mail to reischke@stud.uni-heidelberg.de.

[^1]:    ${ }^{1}$ Note that the only remaining option $\partial^{\mu} \partial_{\mu} \phi$ is a total derivative and will therefore not alter the equations of motion under the usual assumptions on the boundary terms.

[^2]:    ${ }^{2}$ Emmy Noether, 1882-1935.

[^3]:    ${ }^{3}$ To date, the only known fundamental theory that meets this requirement including gravity is string theory.

[^4]:    ${ }^{4}$ In fact, the problem is even more severe as becomes apparent in the Wilsonian approach to be discussed in QFT2. For more information see e.g. the review arXiv:1309.4133 by Cliff Burgess.

[^5]:    ${ }^{1}$ This is simply the statement that the momentum operator is conserved - cf. (1.137).

[^6]:    ${ }^{2}$ Whenever we do not know what to do, and Fourier transformation is not the answer, we insert a $\mathbb{1}$.

[^7]:    ${ }^{3}$ If this were to hold, then the field theory would be free: Indeed the assumption that $\phi(x) \dot{\phi}(y) \rightarrow Z \phi_{\text {in }}(x) \dot{\phi}_{\text {in }}(y)$ and similarly for $\dot{\phi}(y) \phi(x)$ implies that $Z=1$ by exploiting the commutation relations for the fields.

[^8]:    ${ }^{4}$ If we allow also for time derivative terms in the interactions, we should be writing here and in the sequel $i\left[H_{0}, \phi\left(t_{0}, \vec{x}\right)\right]=\Pi_{0}\left(t_{0}, \vec{x}\right)$ with $\Pi_{0}\left(t_{0}, \vec{x}\right) \neq \Pi(t, \vec{x})$. It can be checked that this does not alter the conclusions.

[^9]:    ${ }^{5}$ You must not confuse time ordering $T$ and time $T$ over the next pages.

[^10]:    ${ }^{6}$ Alternatively, the argument can be phrased as follows: Renormalise your theory such that $E_{\Omega}=0$, but $E_{0} \neq 0$. Then it is clear that only the term involving $|\Omega\rangle$ survives unsupressed. The following equations must then be adjusted, but the final result is the same.

[^11]:    ${ }^{7}$ See also the remark right after equ. (2.114).

[^12]:    ${ }^{8}$ The elimination of the $z_{i}$-integration works this way only if the vertex is connected to at least one other point, either internal or external. This is, for example, not the case for the second diagram in Fig. 2.3. More generally, one overall factor of $\int d^{4} z_{i}=\operatorname{Vol}_{\mathbb{R}^{1,3}}=(2 \pi)^{4} \delta^{(4)}(0)$ remains for diagrams not connected to any of the external points. We will see in the next subsection how to deal with such 'disconnected diagrams'.
    ${ }^{9}$ Careful: By momentum space Feynman rules we do not mean that we compute the Fourier transform of the correlator, but rather that we give an equivalent set of rules for the computation of the correlator where the integral over the vertex positions has been performed explicitly.

[^13]:    ${ }^{10}$ In fact, due to renormalisation of the other operators order by order in $\lambda$ this is a bit of an oversimplification, but we will come to this in more detail later in the course.

[^14]:    ${ }^{11}$ Here we are assuming that all particles of the beam hitting the target get to interact with the target particles with equal probability. For simplicity we are also assuming that the particle density varies only in the directions $x^{1}, x^{2}$ transverse to the beam.

[^15]:    ${ }^{12}$ For a more formal, but also considerably more complicated proof see Peskin-Schröder p. 102 - 108 .

[^16]:    ${ }^{1}$ The concept of Lie groups, Lie algebras and their representations has been discussed in detail in the course on Quantum Mechanics. A good book in the present context is: Urbantke, Sexl: Relativity, Groups, Particles. As a quick reminder of the main points see e.g. http://www.thphys.uni-heidelberg.de/\$\sim\$weigand/Skript-QM2011/skript.pdf.
    ${ }^{2}$ In short, a representation $R$ of a group $G$ is a group homomorphism from $G$ to $\operatorname{Aut}(V)$ for some vector space $V$.

[^17]:    ${ }^{3}$ For further reading see e.g. Urbantke, Sexl: Relativity, Groups, Particles.

[^18]:    ${ }^{4}$ More generally, scalars and spinors contribute with opposite signs in loops and theories with supersymmetry, i.e. with an equal number of bosonic and fermionic degrees of freedom, therefore have a chance to exhibit better UV properties.

[^19]:    ${ }^{5}$ For a general proof see Weinberg $(1,5.7)$.

[^20]:    ${ }^{1}$ See Itzykson/Zuber, p. 132 for a proof.

[^21]:    ${ }^{2}$ See Weinberg I, 7.5 for details. The main difference to the massless theory is that $\Pi_{0} \equiv 0$ now poses no problems because, unlike in the massless case, we can solve $A^{0}$ for the spacelike degrees of freedom and simply proceed with the quantisation of $\left(A^{i}, \Pi_{i}\right)$. In more sophisticated terms, the system is amenable to quantisation with Dirac constraints, see again Weinberg.

[^22]:    ${ }^{3}$ Cf. Itzykson, Zuber, p. 136 ff . for details. Note that a priori this action does include spin-0 components in agreement with the above claim that the Proca action is the most general action describing spin-1 degrees of freedom only.

[^23]:    ${ }^{4}$ For instance, sometimes it is useful to consider $\epsilon_{\mu}^{ \pm}(p)= \pm \frac{1}{\sqrt{2}}\left(\epsilon_{\mu}(p, 1) \pm i \epsilon_{\mu}(p, 2)\right)$ to describe photon states of circular polarisation, which coincide with helicity $\pm 1$ states).

[^24]:    ${ }^{5}$ A Yukawa-type deviation of the electromagnetic potential would be directly measurable and constrains the mass of the photon to be below $10^{-14} \mathrm{eV}$, see E. R. Williams, J. E. Faller, and H. A. Hill, Phys. Rev. Lett. 26, 721-724 (1971). In addition, a variety of cosmological and astrophysical constraints imply that the photon mass is at best $10^{-18} \mathrm{eV}$, as reviewed e.g. in http://arxiv.org/pdf/0809.1003v5.pdf

[^25]:    ${ }^{6}$ The remaining Standard Model forces mediated by the eight massless gluons of $S U(3)$ gauge theory is also shortranged, but this is because of confinement.
    ${ }^{7}$ Current constraints imply that the graviton mass must be smaller than $10^{-20} \mathrm{eV}$, see http://arxiv.org/pdf/0809.1003v5.pdf.

[^26]:    ${ }^{1}$ Recall that the indices $\mu, v$ are completely unrelated to the spinor indices, c.f. section 3.2 .

[^27]:    ${ }^{2}$ Note that this is only a question for massless vector fields. For massive vector fields the negative norm states are just the photons with polarization $k^{\mu}\left(\right.$ as $\left.k^{2}>0\right)$ and these obviously decouple by the Ward identity.

[^28]:    ${ }^{3}$ Indeed, the fundamental reason for this divergence is because in QFT the fundamental objects are still pointlike, i.e. they have no substructure. In string theory, on the other hand, the fundamental objects do have an intrinsic substructure (as 1-dimensional strings instead of points) and correspondingly this theory is free of UV divergences.

[^29]:    ${ }^{4}$ In fact, this argument applies to the fully resummed propagator rather than to $\Pi_{\mu \nu}(q)$. However, given the relation between both via Dyson resummation it is not hard to see that the Ward identity carries over to $q^{\mu} \Pi_{\mu \nu}(q)=0$ because it must hold order by order in the coupling constant.
    ${ }^{5}$ To avoid confusion, think of computing the diagram for finite cut-off $\Lambda$ or $\epsilon>0$. The statement is that apart from the UV divergence as the cutoff is removed, $i \Pi^{\mu \nu}(q)$ exhibits no analytic pole at $q^{2}=0$.
    ${ }^{6} \mathrm{We}$ will see this via an easy path-integral proof in the course QFT 2.

[^30]:    ${ }^{7}$ However, once we take $e\left(q^{2}=0\right)$ from experiment, QED does predict the logarithmic running of the effective charge as a function of $q^{2}$.

[^31]:    ${ }^{1}$ Richard Feynman, 1948.

[^32]:    ${ }^{2}$ A careful treatment can be found e.g in Barry Simon, Functional integration and quantum physics, Chelsea Publishing 2005.

[^33]:    ${ }^{3}$ For an overview of some recent developments see e.g. the thesis by C. Anderson (www. math.harvard.edu/theses/senior/anderson/anderson.pdf) and references therein.

[^34]:    ${ }^{4}$ For a detailed treatment of wave functionals we refer e.g. to Chapter 10 of Hatfield, QFT of Point Particles and Strings.

[^35]:    ${ }^{5}$ The most famous one is in fact the formulation of General Relativity as the quantum field theory of a spin- 2 particle, where the classical potential is recovered from scattering amplitudes at 1-loop order in perturbation theory.

[^36]:    ${ }^{6}$ Note that the sum starts at $n=1$ so that $W[0]=0$.

[^37]:    ${ }^{7}$ For a proof we refer e.g. to Brown, Chapter 6.5.

[^38]:    ${ }^{8}$ By contrast, we had $\left.e^{\frac{i}{\hbar} W[J]}=\frac{1}{Z[0]} \int \mathcal{D} \phi e^{\frac{i}{\hbar}(S}[\phi]+J \cdot \phi\right)$.

[^39]:    ${ }^{1}$ Due to Bogoliubov, Parasiuk, Hepp, Zimmermann.

[^40]:    ${ }^{2}$ The factor of $\frac{1}{2}$ in front of the integral is a symmetry factor.

[^41]:    ${ }^{3} \mathrm{Cf}$. renormalisation condition 3. in the previous section, which defines $\lambda$ at scale $\mu$.

[^42]:    ${ }^{4}$ The modification to a fermionic theory with dimensionless couplings is straightforward.

[^43]:    ${ }^{5}$ In the less severe case, where $\lambda \rightarrow \infty$ as $\mu \rightarrow \infty$ one might consider it as a matter of taste whether or not one should call the theory well-defined.
    ${ }^{6}$ A numerical analysis via lattice field theory seems to suggest that the Landau pole in both cases is physical and not just an artifact of perturbation theory.

[^44]:    ${ }^{7}$ By this we mean any product of the fundamental fields or their derivatives.

[^45]:    ${ }^{1}$ For details of this technique see e.g. the book by Kugo.

[^46]:    ${ }^{2}$ Indeed we will assume this here. In actuality, this condition is not satisfied because of a residual gauge symmetry. Consequently, there are several gauge equivalent field configurations that satisfy $F(A)=0$. These are called Gribov copies. One must restrict integration to a fundamental region where the gauge fixing is unique. This restriction indeed has physical consequences, e.g. it modifies the form of the propagator in the deep infra-red.

[^47]:    ${ }^{3}$ Due to Becchi, Rouet, Stora and Tyutin.

[^48]:    ${ }^{4}$ In the sequel $S$ denotes the S-matrix, not the BRST transformation operator. We hope this causes no confusion.

