Elements of Classical Field Theory  
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JOSEPH CONLON\textsuperscript{a}  

\textsuperscript{a}Rudolf Peierls Centre for Theoretical Physics  
University of Oxford  
OX1 3PN Oxford, United Kingdom  

Please send corrections to \texttt{joseph.conlon@physics.ox.ac.uk}.

The parts of the script that are colored blue are not covered in the lecture. They will however be discussed (at least in parts) in the exercises.

The parts of the script that are colored red are strictly non-examinable and will not be covered in the exercises. They represent, for those who are interested, a more advanced gloss on the material.

These notes are adapted from the lectures notes of Uli Haisch.

1 Classical Field Theory

We start by discussing various aspects of classical fields. This will lay the groundwork for the subsequent analysis of quantum fields, as well as illustrating the dynamics of fields beyond what we will be able to analyse in the quantum theory.

1.1 Dynamics of Fields

A field is a quantity that can have a different value at every point of a space(-time) $x = (t, \mathbf{x})$. The spatial coordinate $\mathbf{x}$ may parametrise space itself, but also may restrict to the interior (or surface) of a material - it depends on what the field is describing.

The particular difference of fields is that they involve an infinite number of degrees of freedom. While classical particle mechanics deals with a finite number of \textit{generalized coordinates} $q_a(t)$, indexed by a label $a$, in field theory we are interested in the dynamics of fields

$$\phi_a(t, \mathbf{x}),$$

(1.1)

where both $a$ and $\mathbf{x}$ are considered as labels. We are hence dealing with an infinite number of degrees of freedom (dofs), at least one for each point $\mathbf{x}$ in space. Notice that the concept of position has been relegated from a dynamical variable in particle mechanics to a mere label in field theory.
What does every point in space(-time) mean physically? If we take a field defined on spacetime, then this implies that the value of the field at one point, and the value at another point separated by $10^{-25}$ metres, are physically separate dynamical quantities. There are, furthermore, an infinite number of dynamical quantities going down to ever smaller and smaller scales. By writing a field $\phi(x,t)$, and describing its dynamics, we are also talking about the dynamics of this infinite number of quantities defined on truly microscopic scales.

Actually, we know this must ultimately be wrong physics. In condensed matter systems, this notion is going to break down when we get to distances much smaller than a lattice spacing - at sub-atomic separations, it really does not make sense to talk about the value of a field $\phi(x,t)$. Even in particle physics, where we might think that space makes sense down to arbitrarily small distances, any classical picture of space will still fail at small enough distances (once quantum gravity kicks in).

But yet we still talk about fields with infinite numbers of degrees of freedom, and in practice what happens at the very shortest of distances turns out to decouple when we get to longer distances. The field $\phi(x,t)$ does contain short-distance dynamics in it - but it also drops out from the long-distance questions we want to pose. That this happens is true but deeply non-obvious, and arises from the physics of renormalisation.

**Lagrangian and Action**

The dynamics of the fields is governed by the Lagrangian. In all the systems we will study in this course, the Lagrangian is a function of the fields $\phi_a$ and their derivatives $\partial_\mu \phi_a$ and given by

$$L(t) = \int d^3 x \mathcal{L}(\phi_a, \partial_\mu \phi_a),$$

where the official name for $\mathcal{L}$ is Lagrangian density, but is in practice often simply called the Lagrangian. I am now also going to assume that the field is defined on four dimensions, three spatial and one time (this is an assumption - there are many interesting and physically important field theories that are defined on dimensions less than three).

For any time interval $t \in [t_1, t_2]$, the action corresponding to (1.2) reads

$$S = \int_{t_1}^{t_2} dt \int d^3 x \mathcal{L} = \int d^4 x \mathcal{L}.$$

Recall that in classical mechanics $L$ depends only on $q_a$ and $\dot{q}_a$, but not on the second time derivatives of the generalized coordinates. In field theory we similarly restrict to Lagrangians $\mathcal{L}$ depending on $\phi_a$ and $\dot{\phi}_a$. Furthermore, with an eye on Lorentz invariance, we will only consider Lagrangians depending on $\nabla \phi_a$ and not higher derivatives.

Notice that employing $c = \hbar = 1$, i.e., working with natural units, the dimension of the action is $[S] = 0$. With $[1.3]$ and $[d^4 x] = -4$, it follows that the Lagrangian must

\[1\] If there is no (or only little) room for confusion, we will often drop the arguments of functions and write $\phi_a = \phi_a(x)$ etc. to keep the notation short.
necessarily have $|\mathcal{L}| = 4$. Other objects that we will use frequently to construct Lagrangians are derivatives, masses, couplings, and most importantly fields. The dimensions of the former two objects are $[\partial_\mu] = 1$ and $[m] = 1$, while the dimensions of the latter two quantities depend on the specific type of coupling and field one considers. We therefore postpone the discussion of the mass dimension of couplings and fields to the point when we meet the relevant building blocks.

**Principle of Least Action**

As with particles, the dynamical behavior of fields can be determined by the *principle of least action*. This principle states that when a system evolves from one given configuration to another between times $t_1$ and $t_2$ it does so along a “path” in configuration space for which the action is an extremum (usually a minimum) and hence satisfies $\delta S = 0$. This condition leads to equations of motion:

$$
\delta S = \int d^4x \left\{ \frac{\partial \mathcal{L}}{\partial \phi_a} \delta \phi_a + \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_a)} \delta (\partial_\mu \phi_a) \right\} = \int d^4x \left\{ \left[ \frac{\partial \mathcal{L}}{\partial \phi_a} - \partial_\mu \left( \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_a)} \right) \right] \delta \phi_a + \partial_\mu \left( \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_a)} \delta \phi_a \right) \right\} = 0. \tag{1.4}
$$

The last term is a total derivative and vanishes for any $\delta \phi_a$ that decays at spatial infinity and obeys $\delta \phi_a(t_1, \mathbf{x}) = \delta \phi_a(t_2, \mathbf{x}) = 0$. For all such paths, we obtain the Euler-Lagrange *equations of motion* (EOMs) for the fields $\phi_a$, namely

$$
\partial_\mu \left( \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_a)} \right) - \frac{\partial \mathcal{L}}{\partial \phi_a} = 0. \tag{1.5}
$$

When we come to quantise fields, ultimately the most powerful approach is the use of the path integral formalism, from which the principle of least action arises as a classical limit. However, the easier (but less powerful) approach comes from using a Hamiltonian as taught in ordinary quantum mechanics.

**Hamiltonian Formalism**

The Hamiltonian formalism of classical mechanics starts with the Hamiltonian,

$$
H(q_i, p_i, t) = \sum_i p_i \dot{q}_i(q_i, p_i, t) - L \tag{1.6}
$$

where the conjugate momentum

$$
p_i = \frac{\partial L}{\partial \dot{q}_i}. \tag{1.7}
$$

Field theory is ‘just’ the mechanics of systems with $N \to \infty$ degrees of freedom. For the Hamiltonian formalism of field theory, as the degrees of freedom are $\phi_a(x)$, the conjugate momentum variable (*momentum density*) is $\pi^a(x)$,

$$
\pi^a = \frac{\partial \mathcal{L}}{\partial \dot{\phi}_a}. \tag{1.8}
$$
In terms of \( \pi^a, \dot{\phi}_a, \) and \( \mathcal{L} \) the Hamiltonian density is given by

\[
\mathcal{H} = \pi^a \dot{\phi}_a - \mathcal{L},
\]

where, as in classical mechanics, we must eliminate \( \dot{\phi}_a \) in favor of \( \pi^a \) everywhere in \( \mathcal{H} \). The Hamiltonian then simply takes the form

\[
H = \int d^3 x \mathcal{H}.
\]

Note the expression for \( H \) is precisely as in the case of classical mechanics and is to be regarded as a function of \( \phi_a \) and \( \pi_a \) (no explicit inclusion of \( \dot{\phi}_a \)).

1.2 Noether’s Theorem

A field theory is not required to possess any particular symmetries, and perfectly consistent field theories can be constructed which possess no special symmetry properties.

But, interesting field theories that we care about generally do have extra symmetry properties. Familiar examples are spatial translation, spatial rotation, time translation invariance and Lorentz symmetry. Less familiar examples are internal gauge or global symmetries (which may relate different types of particle to each other), chiral symmetries (multiplying fermions fields by complex phases), or supersymmetries (which relate bosons and fermions).

The role of symmetries in field theory is therefore possibly even more important than in particle mechanics. As with particle mechanics, a crucial aspect of symmetries is that they lead to the generation of conserved quantities. We therefore start here by recasting Noether’s theorem within a field theoretic framework.

Currents and Charges

The statement of Noether’s theorem is that every continuous symmetry of the Lagrangian gives rise to a conserved current \( J^\mu(x) \), so that the EOMs (1.5) imply

\[
\partial_\mu J^\mu = 0,
\]

or in components \( dJ^0/dt + \nabla \cdot J = 0 \). To every conserved current there exists also a conserved (global) charge \( Q \), i.e., a physical quantity which stays the same value at all times, defined as

\[
Q = \int_{R^3} d^3 x J^0.
\]

The latter statement is readily shown by taking the time derivative of \( Q \),

\[
\frac{dQ}{dt} = \int_{R^3} d^3 x \frac{dJ^0}{dt} = -\int_{R^3} d^3 x \nabla \cdot J,
\]

\[\text{In quantum theory there can be exceptions to this, called ‘anomalies’, where symmetries of the classical Lagrangian cease are not symmetries of the quantum theory.}\]
which is zero, provided we assume that \( J \) falls off sufficiently rapidly as \(|x| \to \infty\). Notice, however, that the existence of the conserved current \( J \) is much stronger than the existence of the (global) charge \( Q \), because it implies that charge is in fact conserved \textit{locally}. To see this, we define the charge in a finite volume \( V \) by

\[
Q_V = \int_V d^3x J^0.
\]

(1.14)

Repeating the above analysis, we find

\[
\frac{dQ_V}{dt} = - \int_V d^3x \nabla \cdot J = - \oint_S dS \cdot J,
\]

(1.15)

where \( S \) denotes the area bounding \( V \), \( dS \) is a shorthand for \( n dS \) with \( n \) being the outward pointing unit normal vector of the boundary \( S \), and we have used Gauss’ theorem. In physical terms the result means that any charge leaving \( V \) must be accounted for by a flow of the current 3-vector \( J \) out of the volume. This kind of local conservation law of charge holds in any local field theory.

\textbf{Proof of Theorem}

In order to prove Noether’s theorem, we’ll consider \textit{infinitesimal} transformations. This is always possible in the case of a continuous symmetry. We say that \( \delta \phi_a \) is a symmetry of the theory if, under a field transformation

\[
\phi_a \to \phi_a + \delta \phi_a,
\]

(1.16)

the Lagrangian changes by a total derivative

\[
\delta L(\phi_a) = \partial_\mu J^\mu(\phi_a),
\]

(1.17)

for a set of functions \( J^\mu \). We then consider the transformation of \( L \) under an arbitrary change of field \( \delta \phi_a \). Glancing at (1.4) tells us that in this case

\[
\delta L = \left[ \frac{\partial L}{\partial \phi_a} - \partial_\mu \left( \frac{\partial L}{\partial (\partial_\mu \phi_a)} \right) \right] \delta \phi_a + \partial_\mu \left( \frac{\partial L}{\partial (\partial_\mu \phi_a)} \delta \phi_a \right).
\]

(1.18)

When the EOMs are satisfied than the term in square bracket vanishes so that we are simply left with the total derivative term. For a symmetry transformation satisfying (1.17), the relation (1.18) hence takes the form

\[
\partial_\mu J^\mu = \delta L = \partial_\mu \left( \frac{\partial L}{\partial (\partial_\mu \phi_a)} \delta \phi_a \right),
\]

(1.19)

or simply \( \partial_\mu J^\mu = 0 \) with

\[
J^\mu = \frac{\partial L}{\partial (\partial_\mu \phi_a)} \delta \phi_a - J^\mu,
\]

(1.20)

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which completes the proof. Notice that if the Lagrangian is invariant under the infinitesimal transformation $\delta \phi_a$, i.e., $\delta \mathcal{L} = 0$, then $\mathcal{J}^\mu = 0$ and $J^\mu$ contains only the first term on the right-hand side of (1.20).

We stress that our proof only goes through for continuous transformations for which there exists a choice of the transformation parameters resulting in a unit transformation, i.e., no transformation. An example is a Lorentz boost with some velocity $\mathbf{v}$, where for $\mathbf{v} = 0$ the coordinates $x$ remain unchanged. There are examples of discrete symmetry transformations where this does not occur. E.g., a parity transformation $P$ does not have this property, and Noether’s theorem is not applicable then.

**Energy-Momentum Tensor**

Recall that in classic particle mechanics, spatial translation invariance gives rise to the conservation of momentum, while invariance under time translations is responsible for the conservation of energy. What happens in classical field theory? To figure it out, let’s have a look at infinitesimal translations

$$x^{\nu} \rightarrow x^{\nu} - \epsilon^{\nu} \implies \phi_a(x) \rightarrow \phi_a(x + \epsilon) = \phi_a(x) + \epsilon^{\nu} \partial_{\nu} \phi_a(x),$$

(1.21)

where the sign in the field transformation is plus, instead of minus, because we are doing an active, as opposed to passive, transformation. If the Lagrangian does not explicitly depend on $x$ but only through $\phi_a(x)$ (which will always be the case in the Lagrangians discussed in the course), the Lagrangian transforms under the infinitesimal translation as

$$\mathcal{L} \rightarrow \mathcal{L} + \epsilon^{\nu} \partial_\nu \mathcal{L}.$$  

(1.22)

Since the change in $\mathcal{L}$ is a total derivative, we can invoke Noether’s theorem which gives us four conserved currents $T^\mu_\nu = (\mathcal{J}^\mu)_\nu$ one for each of the translations $\epsilon^{\nu}$ ($\nu = 0, 1, 2, 3$). From (1.20) and (1.21) we readily read off the explicit expressions for $T^\mu_\nu$,

$$T^\mu_\nu = \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_a)} \partial_\nu \phi_a - \delta^\mu_\nu \mathcal{L}. $$

(1.23)

This quantity is called the energy-momentum (or stress-energy) tensor. It has dimension $[T^\mu_\nu] = 4$ and satisfies

$$\partial_\mu T^\mu_\nu = 0.$$  

(1.24)

The four “conserved charges” are ($\mu = 0, 1, 2, 3$)

$$P^\mu = \int d^3 x \mathcal{T}^0\mu,$$

(1.25)

Specifically, the “time component” of $P^\mu$ is

$$P^0 = \int d^3 x T^00 = \int d^3 x \left( \pi^a \dot{\phi}_a - \mathcal{L} \right),$$

(1.26)

which (looking at (1.9) and (1.10)) is nothing but the Hamiltonian $H$. We thus conclude that the charge $P^0$ is the total energy of the field configuration, and it is conserved. In fields theory,
energy conservation is thus a pure consequence of time translation symmetry, like it was in particle mechanics. Similarly, we can identify the charges \( P^i \) \((i = 1, 2, 3)\),

\[
P^i = \int d^3 x T^{0i} = - \int d^3 x \pi^a \phi_a, \tag{1.27}
\]
as the momentum components of the field configuration in the three space directions, and they are of course also conserved.

### 1.3 Symmetries and the Lorentz Group

In general, there are many different symmetry that can occur in physics. A symmetry transformation is, generally, an operation you can perform on a system that transforms it back to itself. In physics, symmetries normally refer to symmetries of the underlying equations. Much interesting physics - including both superconductivity and the properties of the Higgs boson - can be understood as examples of where there is a symmetry of the equations that is not a symmetry of the solution.

We are not going to give a formal mathematical presentation of group theory in these lectures (see the previous C6 lecture notes by John Chalker and Andre Lukas for an account). However, one concept of general importance is that of a representation of a symmetry. Representations refer to the distinct ways an abstract symmetry can be concretely realised, by mapping each symmetry transformation \( g_i \) onto an \( n \times n \) matrix \( G_i \) acting on an \( n \)-dimensional vectors. Examples are the 3-dimensional vector representation of spatial rotations or the ten-dimensional symmetric tensor (metric) representation of Lorentz transformations in general relativity.

One extremely important example of a symmetry is that of the Lorentz group. The Lorentz group \( L \) is of fundamental importance for the construction of relativistic field theories, as \( L \) is associated with the symmetry of 4-dimensional space-time. We therefore now consider the form of the Lorentz group in some more detail.

Using the Minkowski metric \( \eta = \text{diag} (1, -1, -1, -1) \), the Lorentz group \( L \) consists of the real \( 4 \times 4 \) matrices \( \Lambda \) that satisfy

\[
\eta = \Lambda^T \eta \Lambda, \tag{1.28}
\]
which written in components \((\mu, \nu, \rho, \sigma = 0, 1, 2, 3)\) takes the form

\[
\eta^{\mu\nu} = \eta^{\sigma\rho} \Lambda^\mu_\rho \Lambda^\nu_\sigma. \tag{1.29}
\]

It is readily seen that the transformations

\[
x^\mu \rightarrow \Lambda^\mu_\nu x^\nu, \tag{1.30}
\]
with \( \Lambda \) satisfying (1.28) leave the distance \( ds^2 \) invariant. Setting for simplicity the speed of light \( c \) to 1, one has

\[
ds^2 = \eta_{\mu\nu} dx^\mu dx^\nu \rightarrow \eta_{\mu\nu} \Lambda^\mu_\rho dx^\rho \Lambda^\nu_\sigma dx^\sigma = \eta_{\rho\sigma} dx^\rho dx^\sigma = ds^2. \tag{1.31}
\]
The Lorentz transformations (LTs) (1.30) therefore preserve distance as measured by the Minkowski metric.
Figure 1.1: The four disconnected components of the Lorentz group - these consist of 'ordinary' Lorentz transformations that may also be multiplied by Parity and Time-reversal.

Let us examine the consequences of (1.28). First, we take its determinant

\[ \det(\eta) = \det(\Lambda^T) \det(\eta) \det(\Lambda), \quad (1.32) \]

from which we deduce that

\[ \det(\Lambda) = \pm 1. \quad (1.33) \]

The case of \( \det(\Lambda) = +1 \) \((-1)\) corresponds to proper (improper) LTs and the associated subgroup is denoted \( L_+ \) \((L_-)\). This implies that parity or space-inversion \( P = \text{diag} (1, -1, -1, -1) \) as well as time-reversal \( T = \text{diag} (-1, 1, 1, 1) \) are improper LTs and as such part of \( L_- \). Second, we look at the component \( \eta^{00} \), in which case one finds from (1.29) the relation

\[ 1 = \eta^{\sigma \rho} \Lambda_0^\sigma \Lambda_0^\rho = (\Lambda_0^0)^2 - \sum_{i=1,2,3} (\Lambda_i^i)^2. \quad (1.34) \]

It follows that

\[ |\Lambda_0^0| \geq 1. \quad (1.35) \]

When \( \Lambda_0^0 \geq 1 \) the LT is said to be orthochronous and part of \( L_+ \), while \( \Lambda_0^0 \leq -1 \) gives a non-orthochronous LT which belongs to \( L_- \). In consequence, the Lorentz group consists out of four classes of LTs as illustrated in Figure 1.1

Let us have a look at some simple example of LTs: an ordinary spatial rotation by the angle \( \theta \) about the \( z \)-axis, and then a boost by \( v < 1 \) along the \( x \)-axis

\[
\Lambda_\mu^\nu = \begin{pmatrix}
1 & 0 & 0 & 0 \\
0 & \cos \theta & -\sin \theta & 0 \\
0 & \sin \theta & \cos \theta & 0 \\
0 & 0 & 0 & 1
\end{pmatrix}, \quad \Lambda_\mu^\nu = \begin{pmatrix}
\gamma & -\gamma v & 0 & 0 \\
-\gamma v & \gamma & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1
\end{pmatrix},
\]

(1.36)

with \( \gamma = (1 - v^2)^{-1/2} \). In fact, any 3-dimensional rotation

\[
\Lambda_\mu^\nu = \begin{pmatrix}
1 & 0 & 0 \\
0 & 0 & O
\end{pmatrix},
\]

(1.37)
leaves $ds^2$ invariant, since one has $O^T O = O O^T = 1_3$ by definition, if $O$ is an element of the 3-dimensional rotation group $SO(3)$.\footnote{The group of 3-dimensional orthogonal matrices is denoted by $O(3)$, while its subgroup of matrices with determinant $+1$ is called the special orthogonal group $SO(3)$. The difference between the two is the parity transformation $x \to -x$, which is in $O(3)$ but not in $SO(3)$.}

Any Lorentz transformation can be decomposed as the product of a rotation, a boost, space-inversion $P$, and time-reversal $T$. Let us concentrate on the continuous transformations. Since there are three rotations and three boosts, one for each space direction, the continuous LTs are described in terms of six parameters. To find the corresponding six generators, i.e., a basis of transformation matrices that describes infinitesimal rotations and boosts, we write

$$\Lambda = 1_4 + iM,$$

where $M$ are purely imaginary $4 \times 4$ matrices. Inserting this linearized LT into the defining relation \eqref{1.28}, implies

$$M = -\eta M^T \eta .$$

This tells us that the generators $M$ must be anti-symmetric in the space-space components, but symmetric in the space-time components. The space of such matrices has indeed dimension six and is spanned by the set $(i = 1, 2, 3)$

$$J_i = \begin{pmatrix} 0 & 0 \\ 0 & M_i \end{pmatrix} ,$$

$$K_1 = \begin{pmatrix} 0 & i & 0 & 0 \\ i & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \end{pmatrix} , \quad K_2 = \begin{pmatrix} 0 & 0 & i \\ 0 & 0 & 0 \\ i & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix} , \quad K_3 = \begin{pmatrix} 0 & 0 & 0 & i \\ 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \\ i & 0 & 0 & 0 \end{pmatrix} .$$

Here $(M_i)_{jk} = -i \epsilon_{ijk}$ with $\epsilon_{ijk}$ the fully anti-symmetric Levi-Civita tensor ($\epsilon_{123} = +1$) are the generators of $SO(3)$. It follows that the matrices $J_i$ ($K_i$) generate rotations (boosts).

It is a matter of simple algebra to work out the commutation relations of the generators $J_i$ and $K_i$ of the Lorentz group. One obtains

$$[J_i, J_j] = i \epsilon_{ijk} J_k , \quad [K_i, K_j] = -i \epsilon_{ijk} J_k , \quad [J_i, K_j] = i \epsilon_{ijk} K_k ,$$

\begin{table}
\centering
\begin{tabular}{|c|c|c|c|}
\hline
$(j_+, j_-)$ & dimension & name & symbol \\
\hline
$(0, 0)$ & 1 & scalar & $\phi$ \\
$\frac{1}{2}, 0$ & 2 & left-handed Weyl spinor & $\chi_L$ \\
$0, \frac{1}{2}$ & 2 & right-handed Weyl spinor & $\chi_R$ \\
$\frac{1}{2}, 0 \oplus 0, \frac{1}{2}$ & 4 & Dirac spinor & $\psi$ \\
$\frac{1}{2}, \frac{1}{2}$ & 4 & vector & $A_\mu$ \\
\hline
\end{tabular}
\caption{Low-dimensional representations of the Lorentz group.}
\end{table}
where \([a, b] = ab - ba\) is the usual commutator. These relations look very similar to the commutation relations of angular momentum or spin that you should know from quantum mechanics (QM). To make the analogy between the Lorentz group and the Lie group \(SU(2)\) even more explicit, we introduce the following new basis of generators

\[
J_i^\pm = \frac{1}{2} (J_i \pm iK_i) .
\]

In terms of these generators, it is straightforward to verify that the Lorentz Lie algebra takes the form

\[
[J_i^\pm, J_j^\pm] = i\epsilon_{ijk} J_k^\pm , \quad [J_i^\pm, J_j^\mp] = 0 .
\]

This means that \(J_i^+\) and \(J_i^-\) independently obey the Lie algebra of \(SU(2)\). By analogy to the spin quantum number, we can therefore introduce a pair \((j_+, j_-)\) of two spins that characterize the possible representations of the Lorentz group. The states within a representation are further distinguished by the eigenvalues of \(J_3^+\) and \(J_3^-\), which can take the values \(m_+ = -j_+, -j_+ + 1, \ldots, j_+ - 1, j_+\) and \(m_- = -j_-, -j_- + 1, \ldots, j_- - 1, j_-\). The dimension, i.e., the number of distinct states in a given representation \((j_+, j_-)\) is hence \((2j_+ + 1)(2j_- + 1)\).

A list of some basic low-dimensional Lorentz-group representations is provided in Figure 1.2. Field theories in Minkowski space usually require Lorentz invariance and, hence, the Lorentz group is of fundamental importance for such theories. Since it is related to the symmetries of space-time it is often also referred as external symmetry of the theory. The classification of Lorentz group representations provides us with objects which transform in a definite way under LTs and thus these objects are the main building blocks of such field theories. In what follows we will not deal with the spinor fields \(\chi_{L,R}\) and \(\psi\) that describe fermions (quarks and leptons). Our focus will be on scalar \(\phi\) and vector \(A_\mu\) fields, which have very simple transformation properties under the Lorentz group:

\[
\phi \rightarrow \phi , \quad A_\mu \rightarrow \Lambda^\nu_\mu A_\nu .
\]

### 1.4 Scalar Field Theories

We now apply the formalism described above to scalar field theories. These are the simplest class of Lorentz invariant field theories, as they only involve fields transforming trivially under the Lorentz group.

#### A Single Real Scalar Field

The Lagrangian density of a real scalar field \(\phi\) is given by

\[
\mathcal{L} = \frac{1}{2} (\partial_\mu \phi)^2 - V(\phi) ,
\]

where the first term is referred to as kinetic energy (or kinetic term), while the second is known as the scalar potential. Notice that \(\phi\) corresponds to the \((j_+, j_-) = (0, 0)\) representation of the Lorentz group. The dimension of the real scalar field is obviously \([\phi] = 1\).

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4Note that if this is expanded, it looks like \(\dot{\phi}^2 - (\nabla \phi)^2 - V(\phi)\). As \(\nabla \phi\) depends on the values of the individual field degrees of freedom, it can also be viewed as a ‘potential’ term, with the Lagrangian then reflecting the particle mechanics form of \(L = T - U\).
The scalar potential can be written as
\[
V = \frac{1}{2} m^2 \phi^2 + \sum_{n=3}^{\infty} \frac{\lambda_n}{n!} \phi^n. \tag{1.46}
\]
Here \(m\) is the mass of \(\phi\) (this will become clear later on) and the coefficients \(\lambda_n\) are called \textit{coupling constants}. Dimensional analysis tells us that
\[
[m] = 1, \quad [\lambda_n] = 4 - n. \tag{1.47}
\]

The coupling terms in \((1.46)\) fall into three different categories. First, dimension-three terms with \([\lambda_3] = 1\). For such terms, we can define a dimensionless parameter \(\lambda_3/E\), where \(E\) has dimension of mass and represents the energy scale of the process of interest. This means that \(\Delta L_3 = \lambda_3 \phi^3/(3!)\) is a small perturbation for high energies, \textit{i.e.}, \(E \gg \lambda_3\), but a big one at low energies, \textit{i.e.}, \(E \ll \lambda_3\). Such terms are called \textit{relevant}, because they become and are most relevant at low energies which, after all, is where most of the physics that we experience lies. In a relativistic quantum field theory (QFT), we have \(E > m\), which means that we can always make these sorts of perturbations small by taking \(\lambda_3 \ll m\). Second, terms of dimension four with \([\lambda_4] = 0\). \textit{E.g.}, \(\Delta L_4 = \lambda_4 \phi^4/(4!)\). Such terms are small if \(\lambda_4 \ll 1\) and are called \textit{marginal}. Third, operators with dimension of higher than four, having \([\lambda_n] < 0\). In this case the appropriate dimensionless parameters is \((\lambda_n E^{n-4})\) and terms \(\Delta L_n = \lambda_n \phi^n/(n!)\) with \(n \geq 5\) are small (large) at low (high) energies. Such contributions are called \textit{irrelevant}, since in daily life, meaning \(E^{n-4} \ll \lambda_n\), these terms do not matter.

As we will see later, it is typically impossible to avoid high-energy processes in QFT. We hence might expect problems with irrelevant terms (or operators) that become important at high energies. Indeed, these operators lead to \textit{non-renormalizable} QFTs in which one cannot make sense of the infinities at arbitrarily high energies. This does not mean that these theories are useless, it just means that they become incomplete at some energy scale and need to be embedded into an appropriate complete theory \textit{aka} an ultraviolet (UV) completion. Let me also add that the above naive assignment of relevant, marginal, and irrelevant operators is not always carved in stone, since quantum corrections can sometimes change the character of an operator.

\textbf{Low-Energy Description}

In typical applications of QFT only the relevant and marginal couplings are important. This is due to the fact that the irrelevant couplings become small at low energies, as we have seen above. In practice this saves us, since instead of considering the infinite number of coupling terms in \((1.46)\), only a handful are actually needed. \textit{E.g.}, in the case of the real scalar field \(\phi\) described earlier, we only have to take into account two operators, namely \(\Delta L_3 = \lambda_3 \phi^3/(3!)\) and \(\Delta L_4 = \lambda_4 \phi^4/(4!)\), in the low-energy limit.

Let us have a closer look at this issue. Suppose that at some day we discover the true superduper theory \textit{aka} the TOE that describes the world at very high energy scales. Whatever this scale is, let’s call it \(\Lambda\). Since it is an energy scale, we obviously have \([\Lambda] = 1\). What we want to understand are the laws of physics at energy scales \(E\) that we can probe directly.
in a laboratory, which given today’s standards, means $E \ll \Lambda$. Let us further suppose that
at energies of order $E$, the laws of physics are described by a real scalar field This scalar field will have some complicated coupling terms\(^{[1.46]}\), where the precise form is dictated by all the stuff that is going on in the TOE. Can we get an idea about the interactions? Well, we can write our dimensionful coupling constants $\lambda_n$ in terms of dimensionless couplings $g_n$, multiplied by a suitable power of the relevant scale $\Lambda$,

$$\lambda_n = \frac{g_n}{\Lambda^{n-4}}.\quad (1.48)$$

The exact values of the dimensionless couplings $g_n$ depend on the details of the TOE\(^{[9]}\) so we have to do some guesswork. Since the couplings $g_n$ are dimensionless, 1 looks like a pretty good and somehow a natural guess. Since we are not completely sure, let’s say $g_n = \mathcal{O}(1)$. This means that in a laboratory with $E \ll \Lambda$ the coupling terms $\Delta L_n = \lambda_n \phi^n / (n!)$ of \([1.46]\) will be suppressed by powers of $(E/\Lambda)^{n-4}$ if $n \geq 5$. Given the CERN Large Hadron Collider (LHC) energy of around $1\text{ TeV}$, this is a suppression by many orders of magnitude. E.g., for $\Lambda = M_P = 10^{16}\text{ TeV}$ corresponding to the Planck mass $M_P$, one has $E/\Lambda = 10^{-16}$. It is this simple argument based on dimensional analysis that ensures that we need to focus only on the first few terms in the interactions, namely those that are relevant and marginal. It also means that if we only have access to low-energy experiments, it is going to be very difficult to figure out the precise nature of the TOE, because its effects are highly diluted except for the relevant and marginal interactions. Some people therefore call the superduper theory that everybody is looking for, not TOE, but TOENAIL, which stands for “theory of everything not accessible in laboratories”.

**Klein-Gordon Equation**

Applying the Euler-Lagrange equations to \([1.45]\) leads to

$$\partial^\mu \partial_\mu \phi + \frac{\partial V}{\partial \phi} = 0.\quad (1.49)$$

This equation describes the dynamical evolution of the scalar field $\phi$. The Laplacian in Minkowski space is sometimes denoted by $\Box$ and $\partial V/\partial \phi$ is commonly written as $V'$. In this notation, the above equation reads $\Box \phi + V' = 0$.

For non-vanishing (interaction) coefficients $\lambda_3, \lambda_4, \ldots$ the equation is non-linear and solutions to \([1.49]\) are difficult to find. We therefore simply ignore them for the time being and consider only the free case. The dynamics of such a non-interacting massive real scalar field is encoded in the famous *Klein-Gordon equation*:

$$\left(\Box + m^2\right) \phi = 0,\quad (1.50)$$

arising from the Lagrangian

$$\mathcal{L} = \frac{1}{2} (\partial_\mu \phi)^2 - \frac{1}{2} m^2 \phi^2.\quad (1.51)$$

\(^{5}\)Of course, we know that this assumption is plain wrong, since the standard model (SM) of particle physics is a non-abelian gauge theory with chiral fermions, but the same argument applies in that case.

\(^{6}\)If we would know the precise structure of the TOE we could, in fact, calculate the couplings $g_n$. 
To find the solutions to the classical Klein-Gordon equation \((1.51)\), we perform a Fourier transform on the field \(\phi\),

\[
\phi(t, x) = \int \frac{d^3p}{(2\pi)^3} e^{i\mathbf{p} \cdot \mathbf{x}} \phi(t, \mathbf{p}).
\]

(1.52)

In momentum space (1.51) simply reads

\[
\left[ \frac{\partial^2}{\partial t^2} + (\mathbf{p}^2 + m^2) \right] \phi(t, \mathbf{p}) = 0,
\]

(1.53)

which tells us that for each value of \(\mathbf{p}\), the Fourier mode \(\phi(t, \mathbf{p})\) obeys the equation of a simple harmonic oscillator with frequency

\[
\omega_p = \sqrt{|\mathbf{p}|^2 + m^2}.
\]

(1.54)

It is then not difficult to see that the most general solution of the classical Klein-Gordon equation is a linear superposition of simple harmonic oscillators:

\[
\phi(x) = \int \frac{d^3p}{(2\pi)^3} \frac{1}{2\omega_p} (a(p)e^{-ipx} + a^*(p)e^{ipx}),
\]

(1.55)

each vibrating at a different frequency with a different amplitude (the appearance of the specific combination of coefficients \(a(p)\) and \(a^*(p)\) is dictated by the reality of the Klein-Gordon field \(\phi\)). The 4-vector \(\mathbf{p}\) in the exponents is understood as \(p_\mu = (\omega_p, \mathbf{p})\).

The decomposition (1.55) sees each individual Fourier mode evolving as a harmonic oscillator. The main puzzling aspect is to understand why the amplitudes for the Fourier modes have been written as \(\frac{a(p)}{2\omega_p}\) rather than just \(a(p)\).

The reason why this is convenient is to notice that the measure \(\frac{d^3p}{(2\pi)^3} \frac{1}{2\omega_p}\) is a Lorentz-invariant quantity. This follows from “reverse engineering”

\[
\int \frac{d^3p}{(2\pi)^3} \frac{1}{2\omega_p} = \int \frac{d^3p}{(2\pi)^3} \frac{1}{2p_0} \bigg|_{p_0=\omega_p} = \int \frac{d^4p}{(2\pi)^4} \delta(p^2 - p^2 - m^2) \bigg|_{p_0>0} = \int \frac{d^4p}{(2\pi)^4} \delta(p^2 - m^2) \bigg|_{p_0>0}.
\]

(1.56)

From the latter result we can also figure out that the Lorentz-invariant delta function for 3-momenta is

\[
2\omega_p \delta^{(3)}(\mathbf{p} - \mathbf{q}),
\]

(1.57)

since

\[
\int \frac{d^3p}{2\omega_p} 2\omega_p \delta^{(3)}(\mathbf{p} - \mathbf{q}) = 1.
\]

(1.58)

The Lorentz invariance of the measure \(\frac{d^3p}{(2\pi)^3} \frac{1}{(2\omega_p)}\) and its consequences, should be kept in mind, because these features will be important at several occasions later on.

The important physical implication is that quantizing the field \(\phi\) is then the same as quantizing a infinite number of harmonic oscillators. A free field \(\phi\) corresponds to an infinite
number of non-interacting harmonic oscillators, and an interacting field $\phi$ corresponds to an
infinite number of interacting harmonic oscillators. You should already know how this is done in
QM and we will see later on how it is done in QFT. As Sidney Coleman once said: “The career
of a young theoretical physicist consists of treating the harmonic oscillator in ever-increasing
levels of abstraction.”.

**Lorentz Invariance**

We want to confirm that a Lorentz Transformation

$$\phi(x) \rightarrow \phi'(x) = \phi(\Lambda^{-1}x) ,$$  \hspace{1cm} (1.59)

leaves the KG action invariant. Notice that the inverse $\Lambda^{-1}$ appears here because we are
dealing with an active transformation, in which the field is truly shifted. To see why this
means that the inverse appears, it will suffice to consider a non-relativistic example such as
a temperature field. Suppose we start with an initial field $\phi(x)$ which has a hotspot at, say,$x = (1, 0, 0)$. Let’s now make a rotation $x \rightarrow O x$ about the z-axis so that the hotspot ends
up at $x = (0, 1, 0)$. If we want to express the new field $\phi'(x)$ in terms of the old field $\phi(x)$,
we have to place ourselves at $x = (0, 1, 0)$ and ask what the old field looked like at the point
$O^{-1}x = (1, 0, 0)$ we came from. This $O^{-1}$ is the origin of the $\Lambda^{-1}$ factor in the argument of
the transformed field in (1.59).

According to the LT (1.59), the transformation of the mass term is $1/2 m^2 \phi^2(x) \rightarrow 1/2 m^2 \phi^2(x')$ with $x' = \Lambda^{-1}x$. The transformation of $\partial_{\mu}\phi$ is

$$\partial_{\mu}\phi(x) \rightarrow_{\text{transformation}} \partial_{\mu}(\phi(x')) = (\Lambda^{-1})^{\nu}_{\mu}(\partial_{\nu}\phi)(x') . \hspace{1cm} (1.60)$$

Here we have used that

$$\partial_{\mu} = \frac{\partial}{\partial x^\mu} = \partial (x^\nu)^{\nu}_{\mu} \partial (x')^{\nu}_{\nu} = (\Lambda^{-1})^{\nu}_{\mu} \partial'_{\nu} . \hspace{1cm} (1.61)$$

and wrote $\partial'_{\nu}\phi(x') = (\partial_{\nu}\phi)(x')$. Below we will also use the notation $\partial_{\mu}\phi(x')$ to denote the
derivative with respect to $x'$ at $x'$. Using the defining property $(\Lambda^{-1})^{\nu}_{\mu}(\Lambda^{-1})^{\sigma}_{\nu} \eta^{\mu\nu} = \eta^{\rho\sigma}$ of
the LTs, we thus find that the derivative term in the Klein-Gordon Lagrangian then transforms as

$$\frac{1}{2} (\partial_{\mu}\phi(x))^2 \rightarrow \frac{1}{2} (\Lambda^{-1})^{\rho}_{\mu}(\partial_{\rho}\phi)(x')(\Lambda^{-1})^{\sigma}_{\nu}(\partial_{\sigma}\phi)(x') \eta^{\mu\nu}$$

$$= \frac{1}{2} (\partial_{\rho}\phi)(x')(\partial_{\sigma}\phi)(x') \eta^{\rho\sigma}$$

$$= \frac{1}{2} (\partial_{\mu}\phi(x'))^2 . \hspace{1cm} (1.62)$$

Putting things together, we see that the action of the Klein-Gordon theory is indeed Lorentz
invariant,

$$S = \int d^4x \mathcal{L}(x) \rightarrow \int d^4x \mathcal{L}(x') = \int d^4x' \mathcal{L}(x') = S . \hspace{1cm} (1.63)$$

\textsuperscript{8}In passive transformations, it is the coordinates that are changed rather than the field.
Notice that changing the integration variables from $d^4x$ to $d^4x'$, in principle introduces an Jacobian factor $\text{det}(\Lambda)$. This factor is, however, equal to 1 for LT connected to the identity, that we are dealing with (remember that $1_4$ is part of the proper, orthochronous or restricted Lorentz group $L^+_4$).

A similar calculation also shows that, as promised, the EOM of the Klein-Gordon field $\phi$, as given in (1.50), are invariant,

$$\left(\partial^2 + m^2\right)\phi(x) \rightarrow \left(\partial'^2 + m^2\right)\phi'(x)$$

$$= \left[\left(\Lambda^{-1}\right)_{\mu}^\nu \partial_{\nu}(\Lambda^{-1})^{\rho\mu} \partial_{\rho} + m^2\right] \phi(x')$$

$$= \left(\eta^{\nu\rho} \partial'_\nu \partial'_\rho + m^2\right) \phi(x') = 0.$$  \hspace{1cm} (1.64)

Here we have used another common notation for the Laplacian, $\partial^\mu \partial_\mu = \partial^2$.

Notice that the above explicit example shows that the Lagrangian formulation of field theory makes it especially easy to discuss Lorentz invariance, since an EOM is automatically Lorentz invariant if it follows from a Lagrangian that is a Lorentz scalar. This is an immediate consequence of the principle of least action. If a LT leaves the Lagrangian unchanged, then an extremum of the action will be transformed to another extremum of the action.

**Conservation Laws**

Let us derive the Hamiltonian density and the energy-momentum tensor for (1.45). The conjugate momentum of $\phi$ is simply $\pi = \dot{\phi}$, and therefore

$$\mathcal{H} = \frac{1}{2} \pi^2 + \frac{1}{2} (\nabla \phi)^2 + V.$$  \hspace{1cm} (1.65)

For the stress tensor we find from (1.23) the expression

$$T^{\mu\nu} = (\partial^\mu \phi)(\partial^\nu \phi) - \frac{1}{2} \eta^{\mu\nu} (\partial_{\rho} \phi)^2 + \eta^{\mu\nu} V.$$  \hspace{1cm} (1.66)

Notice that the energy-momentum tensor of (1.45) immediately comes out symmetric, i.e., $T_{\mu\nu} = T_{\nu\mu}$. This is not always the case. In accordance with our general formula (1.26), we hence have

$$P^0 = \int d^3x \ T^{00} = \int d^3x \left(\frac{1}{2} \pi^2 + \frac{1}{2} (\nabla \phi)^2 + V\right) = \int d^3x \mathcal{H} = H.$$  \hspace{1cm} (1.67)

The energy is hence conserved in our real scalar theory. Analog relations hold for the space components $P^i$, signalling 3-momentum conservation.

In classical particle mechanics, rotational invariance gives rise to conservation of angular momentum. What is the analogy in field theory? What happens to the remaining three LTs, namely the boosts. What conserved quantity do they correspond to? In order to address these questions using Noether's theorem, we first need the infinitesimal form of the LTs

$$\Lambda^\mu_{\nu} = \delta^\mu_{\nu} + \omega^\mu_{\nu},$$  \hspace{1cm} (1.68)
where $\omega_{\mu\nu}$ is infinitesimal. The condition (1.29) for $\Lambda$ to be a LT becomes in infinitesimal form

$$\eta_{\mu\nu} = \eta^{\rho\sigma} (\delta_{\mu}^\rho + \omega_{\mu}^\rho) (\delta_{\nu}^\sigma + \omega_{\nu}^\sigma) = \eta_{\mu\nu} + \omega_{\mu\nu} + \omega_{\nu\mu} + O(\omega^2).$$

(1.69)

This shows that $\omega_{\mu\nu}$ must be an anti-symmetric matrix (a feature that I have stated before without proof),

$$\omega_{\mu\nu} = -\omega_{\nu\mu}. \quad (1.70)$$

Notice that an anti-symmetric $4 \times 4$ matrix has six independent parameters, which agrees with the number of different Lorentz transformations, i.e., three rotations and three boosts.

Applying the infinitesimal LT to our real scalar field $\phi$, we have

$$\phi(x^\lambda) \rightarrow \phi(x^\lambda - \omega^{\lambda}_{\nu} x^\nu) = \phi(x) - \omega^{\lambda}_{\nu} x^\nu \partial_\lambda \phi(x), \quad (1.71)$$

where the minus sign arises from the factor $\Lambda^{-1}$ in (1.59). The variation of the field $\phi$ under an infinitesimal LT is hence given by

$$\delta \phi = -\omega^{\mu}_{\nu} x^\nu \partial_\mu \phi. \quad (1.72)$$

By the same line of reasoning, one shows that the variation of the Lagrangian (as it is also a Lorentz scalar) is

$$\delta \mathcal{L} = -\omega^{\mu}_{\nu} x^\nu \partial_\mu \mathcal{L} = -\partial_\mu \left( \omega^{\mu}_{\nu} x^\nu \mathcal{L} \right),$$

(1.73)

where in the last step we used the fact that $\omega_{\mu\mu} = 0$ due to its anti-symmetry. The Lagrangian changes by a total derivative, so we can apply Noether’s theorem (1.20) with $\mathcal{J}^\mu = -\omega^{\mu}_{\nu} x^\nu \mathcal{L}$ to find the conserved current,

$$J^\mu = -\partial_\mu \left( \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi)} \omega^{\rho}_{\nu} x^\nu \partial_\rho \phi + \omega^{\mu}_{\nu} x^\nu \mathcal{L} \right)$$

$$= -\omega^{\rho}_{\nu} \left[ \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi)} \partial_\rho \phi - \delta^\mu_\rho \mathcal{L} \right] x^\nu = -\omega^{\rho}_{\nu} T^\mu_{\rho} x^\nu. \quad (1.74)$$

Stripping off $\omega^{\rho}_{\nu}$, we obtain six different currents, which we write as

$$(T^\lambda)^{\mu\nu} = x^\mu T^{\lambda\nu} - x^\nu T^{\lambda\mu}. \quad (1.75)$$

These currents satisfy

$$\partial_\lambda (T^\lambda)^{\mu\nu} = 0, \quad (1.76)$$

and give (as usual) rise to six conserved charges. For $\mu, \nu \neq 0$, the LT is a rotation and the three conserved charges give the total angular momentum of the field configuration ($i, j = 1, 2, 3$):

$$Q^{ij} = \int d^3 x \left( x^i T^{0j} - x^j T^{0i} \right). \quad (1.77)$$

What’s about the boosts? In this case, the conserved charges are

$$Q^{0i} = \int d^3 x \left( x^0 T^{0i} - x^i T^{00} \right). \quad (1.78)$$
The fact that these are conserved tells us that

\[ 0 = \frac{dQ^{0i}}{dt} = \int d^3 x T^{0i} + t \int d^3 x \frac{dT^{0i}}{dt} - \frac{d}{dt} \int d^3 x x^i T^{00} \]

(1.79)

Yet, also the momentum \( P^i \) is conserved, i.e., \( dP^i / dt = 0 \), and we conclude that

\[ \frac{d}{dt} \int d^3 x x^i T^{00} = \frac{d}{dt} \int d^3 x x^i \mathcal{H} = \text{const}. \]

(1.80)

This is the statement that the center of energy of the field travels with a constant velocity. In a sense it’s a field theoretical version of Newton’s first law but, rather surprisingly, appearing here as a conservation law. Notice that after restoring the label \( a \) our results for \((\mathcal{I}^\lambda)^{\mu\nu} \) etc. also apply in the case of multicomponent fields.

**Poincaré Invariance**

We now require that a physical system possesses both space-time translation (1.21) and Lorentz symmetry. The symmetry group that includes both transformations is called the Poincaré group. Notice that for any Poincaré-invariant theory the two charge conservation equations (1.24) and (1.76) should hold.

This is only possible if the energy-momentum tensor \( T^{\mu\nu} \) is symmetric. Applying first (1.76) and then (1.24), we have

\[ 0 = \partial_\lambda (\mathcal{I}^\lambda)^{\mu\nu} = \partial_\lambda (x^\mu T^{\lambda\nu} - x^\nu T^{\lambda\mu}) = x^\mu \partial_\lambda T^{\lambda\nu} - x^\nu \partial_\lambda T^{\lambda\mu} = T^{\lambda\nu} \delta_\lambda^\mu - T^{\lambda\mu} \delta_\lambda^\nu = T^{\mu\nu} - T^{\nu\mu}. \]

(1.81)

This general result tells us that the expression of the energy-momentum tensor in any Poincaré-invariant theory can be made symmetric without changing physics. The key to actually do it, lies in making use of the conservation law (1.24) in an appropriate way.

**Discrete Internal Symmetries**

So far we have only imposed external symmetries, such as Lorentz or Poincaré symmetry, on our scalar theory. Let us also have a brief look at internal symmetries. An interesting (discrete) internal symmetry to be considered is a \( \mathbb{Z}_2 \) symmetry which acts as

\[ \phi(x) \rightarrow -\phi(x). \]

(1.82)

This transformation leaves (1.45) invariant if the scalar potential \( V \) contains only terms with an even number of \( \phi \) fields (but as it is not a continuous symmetry, we cannot use Noether’s
Figure 1.3: Shape of the scalar potential for $m^2 \geq 0$ (solid line) and $m^2 < 0$ (dashed line). In the latter case the positions of the minima are at $\phi = \pm v = \pm \sqrt{(-6m^2)/\lambda}$. The value of the potential at the minima is $V(v) = V_0 - (\lambda v^4)/24$.

Theorem to derive a conserved quantity). Restricting ourselves to relevant and marginal couplings, we write

$$ V = V_0 + \frac{1}{2} m^2 \phi^2 + \frac{\lambda}{4!} \phi^4. \quad (1.83) $$

where we have added a constant term $V_0$ to the potential. Note that the inclusion of the constant term $V_0$ leaves the EOM (1.49) unaltered.

Let us now study the simplest type of solutions to the real scalar theory with the potential (1.83). A trivial solution of (1.49) is $\phi(x) = v = \text{const.}$, provided that

$$ \left. \frac{\partial V}{\partial \phi} \right|_{\phi=v} = 0. \quad (1.84) $$

This relation implies that we should look for the extrema of the potential (1.83). In fact, to minimize the energy (1.67) we should focus on the minima of $V$. Such solutions of the classical theory are called vacua (or vacuum if there is a single one). Note that if the quartic coupling $\lambda$ is negative, the potential is not bounded from below and the energy of a constant field configuration tends to minus infinity for large field values. To avoid such an unphysical situation, we simply assume $\lambda > 0$ in what follows.

The shape of the scalar potential $V$ is shown in Figure 1.3. We see that one has to distinguish two cases. First, the case where $m^2 \geq 0$ (solid curve). In this situation there is a single minimum located at the origin $\phi = 0$. This solution is mapped into itself under the transformation (1.82) and we say that the symmetry is unbroken in this vacuum. The second case with $m^2 < 0$ (dashed curve) is more interesting. In this case $\phi = 0$ is a maximum, while there are two minima. These are located at

$$ \phi = \pm v = \pm \sqrt{-6m^2/\lambda}. \quad (1.85) $$
At the minima, the potential takes the value

\[ V(v) = V_0 + \frac{1}{4} m^2 v^2 = V_0 - \frac{\lambda v^4}{24}. \]  

(1.86)

Notice that neither minimum is invariant under the Z_2 symmetry. In fact, the minima are interchanged under (1.82), and we say that the symmetry is \textit{spontaneously broken}. In general, spontaneous symmetry breaking refers to the case where a symmetry of the underlying theory is partially or fully broken by the vacuum solution of the theory. This concept of spontaneous symmetry breaking is of fundamental importance in theoretical physics and we shall return to it later.

**Continuous Internal Symmetries**

Besides discrete symmetries, continuous internal symmetries also play an important role in particle physics. We start our discussion by considering a scalar theory that involves two real scalar fields \( \phi_{1,2} \) and combine these fields into a field vector

\[ \phi = \begin{pmatrix} \phi_1 \\ \phi_2 \end{pmatrix}. \]  

(1.87)

In the following we will be interested in theories that are invariant under SO(2) transformations \textit{(aka internal 2-dimensional rotations)} that take the form

\[ \phi \rightarrow \phi' = O \phi. \]  

(1.88)

Since this transformation is the same at every space-time point, symmetries of this type are called \textit{global}.\footnote{In a \textit{local} transformation, the matrix \( O \) is allowed to differ between every point in space.} The length of \( \phi \) should be SO(2) invariant, which implies that

\[ O^T O = OO^T = 1_2, \]  

(1.89)

and we also require that

\[ \det (O) = +1. \]  

(1.90)

Proper 2-dimensional rotation matrices can be written as

\[ O = e^{i\alpha T}, \]  

(1.91)

where \( \alpha \) is a real parameter (the rotation angle) and \( T \) denotes the \( 2 \times 2 \) matrix that generates the 2-dimensional rotations. A suitable choice of the \textit{generator} reads

\[ T = \begin{pmatrix} 0 & -i \\ i & 0 \end{pmatrix}. \]  

(1.92)

Check that this indeed exponentiates to a standard rotation! Notice that SO(2) is said to be an \textit{abelian} group, since

\[ e^{i\alpha T} e^{i\beta T} = e^{i\beta T} e^{i\alpha T} = e^{i(\alpha + \beta)T}, \]  

(1.93)
which means that the ordering of two successive group transformations (rotations) does not matter.

For 3-dimensional rotations, things are a bit more complicated. Since there is now one independent rotation for each plane in three dimensions, the SO(3) group has \(3(3-1)/2 = 3\) parameters \(\phi^i\) (Euler angles) and three generators \(T_i\) \((i = 1, 2, 3)\). Infinitesimal SO(3) transformations thus act on \(\phi = (\phi_1, \phi_2, \phi_3)^T\) as \(\phi \to \phi' = \phi + i\alpha^iT_i\phi\). From (1.89) and (1.90) it follows that the generator \(T_i\) are hermitian and traceless:

\[
(T_i)^\dagger = T_i, \quad \text{Tr}(T_i) = 0. \tag{1.94}
\]

A possible set of \(3 \times 3\) matrices with these properties is

\[
T_1 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & -i \\ 0 & i & 0 \end{pmatrix}, \quad T_2 = \begin{pmatrix} 0 & 0 & i \\ 0 & 0 & 0 \\ i & 0 & 0 \end{pmatrix}, \quad T_3 = \begin{pmatrix} 0 & -i & 0 \\ i & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}, \tag{1.95}
\]

which can be written in a more compact way as \((T_i)_{jk} = -i\epsilon_{ijk}\ (j, k = 1, 2, 3)\). By a straightforward calculation, one can prove that the latter matrices satisfy the following Lie algebra

\[
[T_i, T_j] = i\epsilon_{ijk}T_k. \tag{1.96}
\]

These relations imply that two 3-dimensional rotations do not in general "commute", which is equivalent to the non-abelian character of SO(3).

Other examples of non-abelian groups are provided by the special unitary groups \(SU(n)\), which consist of all complex \(n \times n\) matrices \(U\) satisfying

\[
U^\dagger U = UU^\dagger = 1_n, \quad \det(U) = +1. \tag{1.97}
\]

The simplest non-trivial example of a special unitary group is \(SU(2)\). It has \(2^2 - 1 = 3\) generators \(\tau_i\) and close to the identity one can write the group transformations as \(U = 1_2 + it^i\tau_i\). From (1.97) it follows that the generators obey \(\tau_i^\dagger = \tau_i\) and \(\text{Tr}(\tau_i) = 0\). A convenient choice of the generators \(\tau_i\) is

\[
\tau_i = \frac{\sigma_i}{2}, \tag{1.98}
\]

where \(\sigma_i\) are the usual Pauli matrices given by

\[
\sigma_1 = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}, \quad \sigma_2 = \begin{pmatrix} 0 & -i \\ i & 0 \end{pmatrix}, \quad \sigma_3 = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}. \tag{1.99}
\]

Recalling that \(\sigma_i\sigma_j = \delta_{ij}1_2 + i\epsilon_{ijk}\sigma_k\), it is easy to show that \(\text{Tr}(\tau_i\tau_j) = \delta_{ij}/2\) and that the \(SU(2)\) generators fulfil the commutation relations

\[
[\tau_i, \tau_j] = i\epsilon_{ijk}\tau_k. \tag{1.100}
\]

Although they lead to the same commutation relations, \(SU(2)\) and \(SO(3)\) are not actually the same groups. They in fact differ in their global properties. \(SU(2)\) is a 'double cover' of
SO(3) - roughly, it contains two copies of SO(3). This is related to the fact that, under $2\pi$ rotations of SO(3) spinors return to minus themselves. This global minus sign means that spinors should be viewed as representations of SU(2).

Interestingly, (1.96) and (1.100) are the same commutation relations. This means that the matrices $\tau_i$ and $T_i$ form two different representations of the SU(2) Lie algebra. The representation that is classified by the $2 \times 2$ matrices (1.98) is called fundamental. It acts on a complex vector of size two and hence the dimension of the representation space is two. The representation where the generators are $3 \times 3$ matrices of the form $(T_i)_{jk} = -i\epsilon_{ijk}$ (as in (1.95)) is instead called the adjoint representation of SU(2). This representation acts on 3-dimensional real vectors and therefore it has the dimension of SU(2) that is equal to three.

A Complex Scalar Field With a U(1) Symmetry

For the further discussion it will prove convenient to arrange the two real scalar fields that we have meet at the beginning of the previous subsection into a single complex Klein-Gordon field defined as

$$\varphi = \frac{1}{\sqrt{2}} (\phi_1 + i\phi_2) .$$

(1.101)

On this complex representation, the transformation (1.88) acts as follows

$$\begin{pmatrix} \phi_1 \\ i\phi_2 \end{pmatrix} \to \begin{pmatrix} \cos \alpha & -\sin \alpha \\ \sin \alpha & \cos \alpha \end{pmatrix} \begin{pmatrix} \phi_1 \\ i\phi_2 \end{pmatrix} \implies \begin{pmatrix} \phi_1 \\ i\phi_2 \end{pmatrix} \to \begin{pmatrix} \cos \alpha & i\sin \alpha \\ i\sin \alpha & \cos \alpha \end{pmatrix} \begin{pmatrix} \phi_1 \\ i\phi_2 \end{pmatrix} .$$

(1.102)

$$\implies \phi_1 + i\phi_2 \to e^{i\alpha} (\phi_1 + i\phi_2) .$$

This means that all terms in the invariant Lagrangians we are looking for should be unchanged under a rephasing of the $\phi$ field, aka a global U(1) transformation,

$$\varphi \to \varphi' = e^{i\alpha} \varphi , \quad \varphi^* \to (\varphi')^* = e^{-i\alpha} \varphi^* .$$

(1.103)

The different transformation properties of $\varphi$ and $\varphi^*$ tell us that these fields carry charge 1 and $-1$, respectively.

Since $\varphi$ and $\varphi^*$ carry a non-zero charge not all terms polynomial in the fields are allowed in the Lagrangian if we want to respect the symmetry. E.g., a term $\varphi^2$ has charge $+2$ and is hence not invariant under (1.103). In general, we can only allow terms with the same number of $\varphi$ and $\varphi^*$ field, so that the sought U(1) invariant Lagrangian takes the form

$$\mathcal{L} = (\partial_\mu \varphi^*)(\partial^\mu \varphi) - V(\varphi^* \varphi) ,$$

(1.104)

with

$$V(\varphi^* \varphi) = V_0 + m^2 \varphi^* \varphi + \frac{\lambda}{4} (\varphi^* \varphi)^2 .$$

(1.105)

Note that it is essential for the invariance of the kinetic term that the group parameter $\alpha$ is space-time independent. For the EOM for $\varphi$, we find from the Euler-Lagrange equation (1.5),

$$\partial^2 \varphi + \frac{\partial V}{\partial \varphi^*} = \partial^2 \varphi + m^2 \varphi + \frac{\lambda}{2} (\varphi^* \varphi) \varphi = 0 .$$

(1.106)
A similar equation holds in the case of $\phi^*$. The conjugate momenta are $\pi = \dot{\phi}^*$ and $\pi^* = \dot{\phi}$ (note that we treat $\phi$ and $\phi^*$ as independent here - this reflects the fact that $\phi$ contains two degrees of freedom within it). The Hamiltonian then reads

$$H = \pi \dot{\phi} + \pi^* \dot{\phi}^* - L = \pi^* \pi + (\nabla \phi^*) \cdot (\nabla \phi) + V(\phi^* \phi).$$

As this theory is invariant under both translations and Lorentz transformations, the theory has conserved stress and angular momentum tensors which can be obtained in complete analogy with the single scalar field case discussed before. We leave the derivation of the explicit expressions for $T^{\mu \nu}$ and $Q^{ij}$ as an useful exercise.

According to Noether’s theorem the presence of the internal symmetry (1.103) leads to a new type of conserved current which we will now derive. The variations of the fields $\phi$ and $\phi^*$ (treated as independent) under the global $U(1)$ symmetry are

$$\delta \phi = i\alpha \phi, \quad \delta \phi^* = -i\alpha \phi^*.$$  

By plugging these variations into the general formula (1.20) and realizing that in our case $J^\mu = 0$ since $\delta L = 0$, the conserved current and charge is readily obtained:

$$J^\mu = i (\phi^* \partial^\mu \phi - \phi \partial^\mu \phi^*), \quad Q = i \int d^3x (\phi^* \dot{\phi} - \phi \dot{\phi}^*) = i \int d^3x (\phi^* \pi - \phi \pi).$$

There is an ambiguity worth noting, when applying Noether’s theorem to find the conserved charge under the transformation (1.103). Obviously, if $Q$ is conserved, then so is every other operator $c_1 Q + c_2$ with $c_{1,2}$ constant numbers. The expression for $Q$ in (1.109) is therefore unique up to a multiplicative and an additive constant. The multiplicative factor essentially denotes the units in which we measure the charge. Notice that we have already used this ambiguity in (1.109) and simply ignored a factor $-\alpha$. Later we will also learn how the ambiguity on the additive constant is removed.

**Spontaneous Symmetry Breaking**

We would now like to discuss the vacua of the theory, i.e., solutions to the EOM (1.106) with $\phi = v = \text{const}$. As for the real scalar case, when $m^2 \geq 0$ there is a single minimum at $\phi = 0$. This solution is left invariant under the transformations (1.103) and hence the $U(1)$ symmetry is unbroken. For $m^2 < 0$ the shape of the potential is illustrated in Figure 1.4. We can see from the depicted red line that in this case there is a whole circle of minima

$$\varphi = v = \frac{v_0}{\sqrt{2}} e^{i\delta}, \quad v_0 = \sqrt{-\frac{4m^2}{\lambda}},$$

where $\delta$ is an arbitrary phase. Clearly, the existence of this one-dimensional space of vacua is not an accident, but originates from the $U(1)$ invariance of the scalar potential (1.105):

$$V(\varphi^* \varphi) \rightarrow V((\varphi')^* \varphi') = V(e^{-ia} \varphi^* e^{ia} \varphi) = V(\varphi^* \varphi).$$

In fact, this symmetry-induced invariance implies that for every minimum $\varphi$ of $V$, $\varphi' = \exp(i\alpha) \varphi$ is also a minimum for arbitrary $\alpha$. We can use this freedom to choose a particular
minimum that breaks the $U(1)$ symmetry spontaneously. Let’s take the minimum on the $\phi_1$ axis, meaning $\varphi = v_0/\sqrt{2}$. Around this point we then can expand the field as

$$\varphi = \frac{1}{\sqrt{2}} (v_0 + \varphi_1 + i\varphi_2), \quad (1.112)$$

where $\varphi_{1,2}$ are understood to be small fluctuations. Inserting the above relation into (1.105), we find

$$V = V_0 + \frac{1}{4} m^2 v_0^2 - m^2 \varphi_1^2 + O(\varphi_1^3, \varphi_1 \varphi_2^2). \quad (1.113)$$

This is an interesting result, because it tells us that while the excitation $\varphi_1$ is massive with mass $\sqrt{2}m$ (remember $m^2 < 0$) the excitation $\varphi_2$ remains massless. Looking at Figure 1.4, this is not a completely unexpected feature. While $\varphi_2$ corresponds to a fluctuation around the circle of minima (a flat direction in the potential), $\varphi_1$ represents an excitation perpendicular to $\varphi_1$ (a direction of the potential with a curvature). The appearance of massless scalars for spontaneously broken global symmetries is a general feature known as Goldstone’s theorem and the corresponding massless scalars are also called Goldstone bosons. We will now study this phenomenon in a more general setting.\(^{10}\)

Goldstone’s Theorem

To generalize the above findings, we consider a field theory involving a set of scalar fields $\phi = (\phi_1, \ldots, \phi_n)^T$ with corresponding potential $V(\phi)$. We assume that the potential has a minimum at $\mathbf{v} = (v_1, \ldots, v_n)^T$ that satisfies

$$\left. \frac{\partial V}{\partial \phi_a} \right|_{\phi = \mathbf{v}} = 0. \quad (1.114)$$

\(^{10}\)Although we are meeting it here as a classical result, Goldstone’s theorem also extends to quantum theories.
After defining $\varphi = \phi - v$, we can Taylor expand the potential around such a minimum in the following way
\[
V = V(v) + \frac{1}{2} M_{ab} \varphi^a \varphi^b + O(\varphi^3),
\]
where $(a, b = 1, \ldots, n)$
\[
M_{ab} = \frac{\partial^2 V}{\partial \phi_a \partial \phi_b} \bigg|_{\phi=v}.
\]
Clearly, the eigenvalues of the mass matrix $M$ are the mass squares of the fields (around the vacuum $v$).

Now let us assume that our scalar theory is invariant under a continuous (global) symmetry group $G$, under which
\[
\phi \rightarrow R(g)\phi,
\]
where $R(g)$ denotes the representation of $\phi$. The invariance of the scalar theory in particular means that the scalar potential is invariant under such a transformation, namely
\[
V(\phi) = V(R^{-1}(g)\phi),
\]
for all $g \in G$. The vacuum will in general not respect the full symmetry group $G$, but will spontaneously break it into a subgroup $H \subset G$, so that
\[
R(g)v = v \text{ for } g \in H, \quad R(g)v \neq v \text{ for } g \notin H.
\]
Now introducing infinitesimal group transformations
\[
R(g) = 1 + it^\alpha T_\alpha,
\]
with generators $T_\alpha$ (in the representation $R$) and small coefficients $t^\alpha$. The set of generators $\{T_\alpha\}$ can be split into two subsets $T_i$ and $T_j$, where $T_i$ denotes the generators of the unbroken subgroup $H$, while $T_j$ are the remaining generators corresponding to the broken part of $G$. By definition, one has
\[
T_i v = 0, \quad T_j v \neq 0.
\]
Now let us see how (1.120) acts on our potential. We have
\[
V(\phi) = V(\phi - it^\alpha T_\alpha \phi) = V(\phi) - i \left( \frac{\partial V}{\partial \phi} \right)^T t^\alpha T_\alpha \phi.
\]
Differentiating this expression again and evaluating it at $\phi = v$, one finds using (1.114) and (1.116), that
\[
MT_\alpha v = 0.
\]
From this relation we see that every broken generator $T_j$ (with $T_j v \neq 0$) leads necessarily to an eigenvector of $M$ with eigenvalue zero. Or in more physical terms, every broken symmetry leads to a massless scalar excitation. We have proven Goldstone’s theorem!
Scalar Field Theory With a $SU(2) \times U(1)$ Symmetry

Let us further illustrate Goldstone’s theorem by an educated example (this is a simplified account of what actually happens in the Higgs sector of the Standard Model). We consider the symmetry $SU(2) \times U(1)$ and a complex scalar field $\Phi$ that is a doublet under $SU(2)$ and has a charge of $1/2$ under the $U(1)$. Our scalar transforms as

$$\Phi \rightarrow e^{i\alpha Y} e^{it\tau_i} \Phi \simeq (1 + i\alpha Y + it\tau_i) \Phi,$$

with the generators

$$Y = \frac{1}{2}, \quad \tau_i = \frac{\sigma_i}{2}.$$

Notice that our notation is on purpose sloppy here, since the 1 on the right-hand side of (1.124) should be in fact the $2 \times 2$ matrix $1_2$. A similar statement also applies to the generator $Y$. Such a sloppiness in notation is common in textbooks, and it is good to get used to it early.

We can write an invariant Lagrangian as

$$L = (\partial_\mu \Phi)^\dagger (\partial^\mu \Phi) - V(\Phi^\dagger \Phi)$$

with

$$V(\Phi^\dagger \Phi) = V_0 + \mu^2 \Phi^\dagger \Phi + \lambda (\Phi^\dagger \Phi)^2.$$

Provided that $\lambda > 0$ and $\mu^2 < 0$ (note we use $\mu^2$ rather than $m^2$ for the squared mass term), the scalar potential (1.127) is bounded from below and minimized for

$$\Phi^\dagger \Phi = v^2_0 = -\frac{\mu^2}{2\lambda}.$$

Using (1.124) allows us to choose a particular simple vacuum configuration. Without loss of generality, we take

$$\Phi = v = \begin{pmatrix} 0 \\ v_0 \end{pmatrix}.$$

For this choice, it is easy to verify that

$$\tau_1 v \neq 0, \quad \tau_2 v \neq 0, \quad (\tau_3 - Y) v \neq 0,$$

but

$$(\tau_3 + Y) v = 0.$$

Hence three of the four generators of $SU(2) \times U(1)$ are broken, which corresponds to the breaking pattern

$$SU(2) \times U(1) \rightarrow U(1),$$

where the $U(1)$ on the right-hand side is the diagonal subgroup of $SU(2) \times U(1)$ corresponding to the generator $Q = \tau_3 + Y$.

Although this example appears slightly contrived, its relevance stems from the fact that the breaking pattern (1.132) is precisely that of the electroweak gauge group of the SM of
particle physics: the SM is in fact based on $SU(3)_c \times SU(2)_L \times U(1)_Y \rightarrow SU(3)_c \times U(1)_{em}$, where the subscript $c$ refers to color, $L$ indicates that the $SU(2)$ only acts on left-handed fields, $Y$ represents hypercharge, and $U(1)_{em}$ corresponds to the unbroken electromagnetism associated to the generator $Q$, i.e., the electric charge.

However, there is an important difference. While the above theory predicts three massless Goldstone bosons, in Nature we only observe a single weakly-interacting massless state, i.e., the photon. This difference is explained by the fact that the electroweak symmetry in the SM is realized as a local or gauge(d) symmetry, while in our simple example we dealt with a global symmetry. It turns out that in the case of the spontaneously broken electroweak gauge symmetry our three massless Goldstone bosons are absorbed (or eaten) by three vector bosons (the electrical neutral $Z$ boson and the charged $W^\pm$ bosons) that receive their mass from the breaking mechanism. This is the famous Higgs mechanism.

The residual massive scalar mode becomes massive is then the (physical) massive scalar state the Higgs. After many years of searching, the Higgs boson was discovered in 2013 when the ATLAS and the CMS collaborations (situated at the CERN LHC in Geneva) announced the existence of a new bosonic state with a mass of around 125 GeV. Subsequent analyses have shown that this particle behaves precisely as expected for a Standard Model Higgs particle. This discovery marks a turning point in the history of elementary-particle physics: the almost 50 year-long hunt for the came to a successful end, and a no-lose era ended (in the sense that the Higgs was a guaranteed discovery for a collider that was energetic enough). The Nobel prize was then awarded to Francois Englert and Peter W. Higgs for their theoretical work on the mechanism of electroweak symmetry breaking.

However, in order to fully understand the physics of the Higgs mechanism, we are still missing an important ingredient, namely vector fields. These 4-component fields will be discussed now.

### 1.5 Vector Field Theories

In this section we focus on the basic features of another important ingredient of realistic (quantum) field theories, i.e., vector fields. These fields carry spin one and are called vector fields, since they transform like a vector under the Lorentz group.

#### Maxwell’s Theory Of Electromagnetism

As another simple application of the formalism we have developed, let us try to derive Maxwell’s equations using the field theory formulation. In terms of the electric and magnetic fields $E$ and $B$ and the charge density $\rho$ and 3-vector current $j$, these equations take the well-known form

$$\nabla \cdot B = 0,$$  \hspace{1cm} (1.133)

$$\nabla \times E + \frac{\partial B}{\partial t} = 0,$$  \hspace{1cm} (1.134)

$$\nabla \cdot E = \rho,$$  \hspace{1cm} (1.135)
The $E$ and $B$ fields are spatial 3-vectors and can be expressed in terms of the components of the 4-vector field $A^\mu = (\phi, A)$ by

$$E = -\nabla \phi - \frac{\partial A}{\partial t}, \quad B = \nabla \times A.$$  

(1.137)

This definition ensures that the first two homogeneous Maxwell equations (1.133) and (1.134) are automatically satisfied,

$$\nabla \cdot (\nabla \times A) = \epsilon_{ijk} \partial_i \partial_j A_k = 0,$$  

(1.138)

$$\nabla \times (-\nabla \phi - \frac{\partial A}{\partial t}) + \frac{\partial}{\partial t} (\nabla \times A) = -\nabla \times (\nabla \phi) = -\epsilon_{ijk} \partial_j \partial_k \phi = 0.$$  

(1.139)

Here the indices $i, j, k = 1, 2, 3$ are summed over if they appear twice.

The remaining two inhomogeneous Maxwell equations (1.135) and (1.136) follow from the Lagrangian

$$L = -\frac{1}{4} F_{\mu\nu} F^{\mu\nu} - A_\mu J^\mu,$$  

(1.140)

with $J^\mu = (\rho, j)$. From the rules presented before, we gather that the dimension of the vector field and current is $[A_\mu] = 1$ and $[J^\mu] = 3$, respectively. The funny minus sign of the first term on the right-hand side is required to ensure that the kinetic term $1/2 \dot{A}_0^2$ is positive using the Minkowski metric. Notice also that the Lagrangian (1.140) has no kinetic term $1/2 \dot{A}_0^2$ and hence $A_0$ is not dynamical. Why this is and necessarily has to be the case will become clear pretty soon.

Let us now evaluate the equations of motion and so confirm that the statement made before (1.140) is indeed correct. We first evaluate

$$\frac{\partial L}{\partial A_\nu} = -J^\nu, \quad \frac{\partial L}{\partial (\dot{A}_\mu A_\nu)} = -\partial^\nu A^\mu + \partial^\nu A^\mu,$$  

(1.141)

from which we derive the EOMs,

$$0 = \partial_\mu \left( \frac{\partial L}{\partial (\dot{A}_\mu A_\nu)} \right) - \frac{\partial L}{\partial A_\nu} = \partial_\mu (-\partial^\nu A^\mu + \partial^\nu A^\mu) + J^\nu.$$  

(1.142)

Introducing now the field-strength tensor

$$F_{\mu\nu} = \partial_\mu A_\nu - \partial_\nu A_\mu,$$  

(1.143)

we can write (1.140) and (1.142) quite compact,

$$L = -\frac{1}{4} F_{\mu\nu} F^{\mu\nu} - J_\mu A^\mu.$$  

(1.144)

$$\partial_\mu F^{\mu\nu} = J^\nu.$$  

(1.145)
Does this look familiar? I hope so. Notice that we have \[ F_{\mu\nu} = 2 \] and that the field-strength tensor satisfies the Bianchi identity:

\[
\partial_\rho F_{\mu\nu} + \partial_\mu F_{\nu\rho} + \partial_\nu F_{\rho\mu} = 0.
\] (1.146)

In order to see that (1.145) indeed captures the physics of (1.135) and (1.136), we compute the components of \( F^{\mu\nu} \). We find

\[
F^{0i} = -F^{i0} = \partial^0 A^i - \partial^i A^0 = \left( \nabla \phi + \frac{\partial A}{\partial t} \right)^i = -E^i,
\] (1.147)

\[
F^{ij} = -F^{ji} = \partial^i A^j - \partial^j A^i = -\epsilon^{ijk} B^k,
\]

while all other components are zero. With this in hand, we then obtain from \( \partial_\mu F^{\mu 0} = \rho \) and \( \partial_\mu F^{\mu 1} = j^1 \),

\[
\partial_\mu F^{\mu 0} = \partial_0 F^{00} + \partial_i F^{i0} = \nabla \cdot E = \rho,
\]

\[
\partial_\mu F^{\mu 1} = \partial_0 F^{01} + \partial_i F^{i1} = -\frac{\partial E^1}{\partial t} + \frac{\partial B^3}{\partial x_2} - \frac{\partial B^2}{\partial x_3} = \left( \nabla \times B - \frac{\partial E}{\partial t} \right)^1 = j^1.
\] (1.148)

Similar relations hold for the remaining components \( i = 2, 3 \). Taken together this proves the second inhomogeneous Maxwell equation (1.136).

Let me also derive the energy-momentum tensor \( T^{\mu\nu} \) of electrodynamics, ignoring for the moment the source term \( J_\mu A^\mu \). Using (1.141) one finds

\[
T^{\mu\nu} = -F^{\mu\rho} \partial_\nu A_\rho + \frac{1}{4} \eta^{\mu\nu} F_{\rho\sigma} F^{\rho\sigma}.
\] (1.149)

Notice that the first term in (1.149) is not symmetric, which implies that \( T^{\mu\nu} \neq T^{\nu\mu} \). In fact, this is not really surprising since the definition of the energy-momentum tensor (1.23) does not exhibit an explicit symmetry in the indices \( \mu \) and \( \nu \). Nevertheless, there is always a way to massage the energy-momentum tensor of any theory into a symmetric form.\(^{11}\) To learn how this can be done in the case under consideration is the objective of a homework problem.

### Physical Degrees Of Freedom

Under LTs, \( A_\mu \) transforms as a vector, \( i.e., A_\mu \rightarrow \Lambda_\mu^\nu A_\nu \), and the same applies of course also to the current \( J_\mu \). The field-strength tensor \( F_{\mu\nu} \), on the other hand, carries two indices and hence transforms as a tensor, \( i.e., F_{\mu\nu} \rightarrow \Lambda_\mu^\rho \Lambda_\nu^\sigma F_{\rho\sigma} \). Equipped with these transformation properties it is an easy exercise to show that the Lagrangian (1.144) is Lorentz invariant (at the end it is sufficient to observe that in the expression for \( \mathcal{L} \) all indices are contracted, so that (1.144) is a Lorentz scalar). The explicit calculation is left as an exercise to the interested reader.

\(^{11}\)One (but not the only) reason that you might want to have a symmetric energy-momentum tensor \( T^{\mu\nu} \) is to make contact with general relativity, since such an object sits on the right-hand side of Einstein’s field equations.
After (1.140) we have already mentioned that not all components of $A_\mu$ are dynamical, which signals that $A_\mu$ contains unphysical dofs. In fact, this is not really a surprise because the physical fields $E$ and $B$ come from $F_{\mu\nu}$ not $A_\mu$ (see (1.147)). Formally, this is expressed by the fact that a gauge transformation

$$A_\mu(x) \rightarrow A_\mu(x) + \partial_\mu f(x),$$

with arbitrary function $f(x)$ leaves $F_{\mu\nu}$ invariant:

$$F_{\mu\nu} = \partial_\mu A_\nu - \partial_\nu A_\mu \rightarrow \partial_\mu (A_\nu + \partial_\nu f) - \partial_\nu (A_\mu + \partial_\mu f)$$

$$= \partial_\mu A_\nu - \partial_\nu A_\mu + \partial_\mu \partial_\nu f - \partial_\nu \partial_\mu f = F_{\mu\nu}. \quad (1.151)$$

A gauge transformation should be thought of relating two physically identical configurations for $A_\mu(x)$.

This is an important result, because it gives us a recipe on how to “derive” the covariant formulation of electrodynamics without ever talking about Maxwell’s equations. The idea is simply to find the most general Lagrangian (including terms up to a certain mass dimension) that is Lorentz invariant and unchanged under (1.150). Notice that gauge invariance in particular implies that $\mathcal{L}$ should only depend on $A_\mu$ through $F_{\mu\nu}$. Up to dimension four this leaves basically only one term (in the case of a vanishing source term), namely $\mathcal{L} = -1/4 F_{\mu\nu} F^{\mu\nu}$. The associated EOM is of course $\partial_\nu F^{\mu\nu} = 0$, or written in terms of $A_\mu$,

$$\partial^2 A_\mu - \partial_\mu \partial_\nu A^\nu = 0. \quad (1.152)$$

In fact, this equation can be further simplified by exploiting the gauge symmetry (1.150) by imposing a gauge condition on $A_\mu$ through a suitable choice of gauge parameter $f$. There are various possibilities for such a gauge choice and here we consider the Lorentz (sic) gauge defined by

$$\partial_\mu A^\mu = 0. \quad (1.153)$$

This gauge has the salient benefit that it is covariant (in contrast to other gauges such as temporal gauge $\dot{A}_0 = 0$ or Coulomb gauge $\nabla \cdot A = 0$) and it simplifies (1.152) to

$$\partial^2 A_\mu = 0. \quad (1.154)$$

It is however important to realize that (1.153) does not fix the gauge entirely but leaves a residual gauge freedom, since all choices of $f$ with

$$\partial^2 f = 0, \quad (1.155)$$

are equivalent in the sense that they all lead to (1.154).

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12The term $F_{\mu\nu} F^{\mu\nu}$ with $\tilde{F}^{\mu\nu} = \epsilon^{\mu\nu\rho\sigma} F_{\rho\sigma}$ denoting the dual field-strength tensor also meets all requirements. It can however be written as a total derivative $4 \partial_\mu (\epsilon^{\mu\nu\rho\sigma} A_\nu \partial_\rho A_\sigma)$ and therefore does not contribute to the classical EOMs. We hence ignore such a term.
The solutions of $\partial^2 A_\mu = 0$ can be quickly obtained by noticing that the latter equation is nothing but a Klein-Gordon equation for a massless vector field. Adding vector indices to (1.55), we arrive at

$$A_\mu(x) = \int \frac{d^3 p}{(2\pi)^3} \frac{1}{2\omega_p} (a_\mu(p)e^{-ipx} + a^*_\mu(p)e^{ipx}) ,$$

where $\omega_p = |p|$ and $p^\mu = (\omega_p, p)$. After a Fourier transform, the Lorenz gauge condition (1.153) can be written as

$$p^\mu a_\mu(p) = 0.$$

This gauge condition can be used to show that there are only two physical degrees of freedom in the electromagnetic field - the polarisation modes transverse to the direction of propagation. This result should be familiar to you from the study of sunglasses, and we now see how it arises here.

To exploit this constraint one conveniently introduces a set of polarization vectors $\epsilon^{(\alpha)}_\mu(p)$ with $\alpha = 0, 1, 2, 3$ and the following properties. The vectors $\epsilon^{(1)}_\mu$ and $\epsilon^{(2)}_\mu$ are orthogonal to both $p$ and a vector $n$ with $n^2 = 1$ and $n_0 > 0$. They furthermore are chosen such that they satisfy $\epsilon^{(\alpha)} \cdot \epsilon^{(\beta)} = -\delta^{\alpha\beta}$ for $\alpha, \beta = 1, 2$. The polarization vector $\epsilon^{(3)}_\mu$ is taken to be in the $p-n$ plane, orthogonal to $n$ and normalized to -1, i.e., $n \cdot \epsilon^{(3)} = 0$ and $(\epsilon^{(3)})^2 = -1$. Finally, one defines $\epsilon^{(0)} = n$. With these conventions one has an orthogonal set of vectors satisfying

$$\epsilon^{(\alpha)} \cdot \epsilon^{(\beta)} = \eta^{\alpha\beta},$$

for $\alpha, \beta = 0, 1, 2, 3$. This basis can be used to write

$$a_\mu(p) = \sum_{\alpha=0}^3 a^{(\alpha)}(p) \epsilon^{(\alpha)}_\mu(p).$$

In fact, the basic idea behind all this “mumbo-jumbo” is that this specific choice of basis of polarization vectors, allows to easily separate the directions transversal to $p$ (corresponding to $\epsilon^{(1)}_\mu$ and $\epsilon^{(2)}_\mu$) from the other two directions (corresponding to $\epsilon^{(0)}_\mu$ and $\epsilon^{(3)}_\mu$) that are longitudinal to $p$. As an example, if we choose a spatial momentum $p$ pointing in the $z$-directions, the components of above vectors are explicitly given by

$$\epsilon^{(\alpha)}_\mu = \delta^{\alpha\mu}.$$

In the general case, one has instead

$$\epsilon^{(3)}_\mu = \frac{p_\mu}{p_0} - n_\mu.$$

Clearly,

$$n \cdot \epsilon^{(3)} = \frac{np}{p_0} - n^2 = \frac{np}{p_0} - 1 = 0 \implies \frac{np}{p_0} = 1,$$

and thus

$$(\epsilon^{(3)})^2 = \frac{p_0}{p_0} + n^2 - 2 \frac{np}{p_0} = 1 - 2 \frac{np}{p_0} = -1,$$
where we have used that our vector field is massless, i.e., $p^2 = 0$. So (1.161) is indeed the sought solution. Notice also that

$$\epsilon^0 + \epsilon^3 = \frac{p_\mu}{p_0}.$$

(1.164)

We now return to (1.157) and evaluate this sum inserting the expansion (1.159). By definition the transversal components $\epsilon^{(1)}$ and $\epsilon^{(2)}$ do not contribute and we find

$$p^\mu a_\mu(p) = (p \cdot \epsilon^0) a^{(0)}(p) + (p \cdot \epsilon^3) a^{(3)}(p) = p_0 \left( a^{(0)}(p) - a^{(3)}(p) \right) = 0.$$

(1.165)

Here we have used that

$$p \cdot \epsilon^0 = -p \cdot \epsilon^3 = p_0.$$

(1.166)

Note that the first equality follows from (1.164) after contracting it with $p^\mu$, while the second equality is a consequence of (1.162). From (1.165) it follows that $a^{(0)}(p) = a^{(3)}(p)$. The expansion (1.159) can thus be written as

$$a_\mu(p) = a^{(3)}(p) \frac{p_\mu}{p_0} + \sum_{\alpha=1}^{2} a^{(\alpha)}(p) \epsilon^{(\alpha)}(p),$$

(1.167)

and we are left with two transversal modes and a longitudinal one along the $p$ direction (remember that after (1.140) we already mentioned the fact that the time component $A_0$ is not dynamical, which reduces the number of independent dofs in $A_\mu$ from four to three). This is still one more dof than a physical massless field, such as the photon, ought to have.

The trick to get rid of the remaining longitudinal component is to make use of the residual gauge freedom (1.155) and to “gauge away” the term proportional to $a^{(3)}(p)$ in (1.167). Let see how this works. Given that $f$ also fulfils a Klein-Gordon equation, the most general decomposition obviously reads

$$f(x) = \int \frac{d^3 p}{(2\pi)^3} \frac{1}{2\omega_p} \left( \xi(p) e^{-ipx} + \xi^*(p) e^{ipx} \right).$$

(1.168)

A gauge transformation (1.150) with such a $f$ corresponds to

$$a_\mu(p) \rightarrow a'_\mu(p) = a_\mu(p) - ip_\mu \xi(p).$$

(1.169)

But such a shift is exactly what is needed to remove the remaining longitudinal component from (1.167) and we arrive at the correct physical answer. Needless to say that the reduction from apparent four dofs to actually two is intimately related to the gauge invariance of (quantum) electrodynamics.\(^{13}\)

A Massive Vector Field

What happens if we try to give our vector field $A_\mu$ a mass? We can try and add a term

$$\frac{1}{2} m^2 A_\mu A^\mu,$$

(1.170)

\(^{13}\)Although we have presented a classical argument, an analogous argument also holds in the quantum theory.
to our Lagrangian $\mathcal{L} = -1/4 F_{\mu\nu} F^{\mu\nu}$, but immediately realize that such a term is not allowed by gauge invariance (1.150). This is of course stupid, but let us ignore this unwanted feature for a moment. The EOMs for the massive vector field would not be $\partial_\mu F^{\mu\nu} = 0$, but

$$\partial_\mu F^{\mu\nu} + m^2 A^\nu = 0.$$  \hfill (1.171)

Applying $\partial_\nu$ to this relation we conclude that $\partial_\nu A^\nu = 0$, since $\partial_\nu \partial_\mu F^{\mu\nu}$ is trivially zero. We can hence split the latter equation into two

$$(\partial^2 + m^2) A_\mu = 0, \quad \partial_\mu A^\mu = 0.$$  \hfill (1.172)

The first equation is a massive Klein-Gordon equation for $A_\mu$ with the general solution (1.156) but $\omega_p = (|p|^2 + m^2)^{1/2}$ instead of $\omega_p = |p|$. In order to satisfy the second equation we need to impose the condition (1.157), which reduces the number of dofs from four to three. Can we scrap another dof? Not this time - the mass term (1.170) breaks gauge invariance and thus we do not have the freedom to use (1.155) and gauge it away. This means that a massive vector field has three physical degrees of freedom, one more than a massless one.

Current Sources

In the discussion of the last two subsections we have always ignored the possible presence of external sources by simply setting $J_\mu = 0$ by hand. If we now include a source term, the relevant Lagrangian leading to Maxwell's theory has already been given in (1.144). By applying (1.150), we see that a gauge transformation of such a term introduces an additional piece into the action:

$$S \rightarrow S - \int d^4 x (\partial_\mu f) J^\mu = S + \int d^4 x f (\partial_\mu J^\mu).$$  \hfill (1.173)

As usual we have used integration by parts to arrive at the final answer. This result tells us that the apparent breaking of gauge invariance by the term $J_\mu A^\mu$ can be avoided, provided that

$$\partial_\mu J^\mu = 0.$$  \hfill (1.174)

In other words, the vector field $A_\mu$ needs to couple to a conserved current $J_\mu$. In a fundamental theory, we expect the current $J_\mu$ to itself arise from fields, rather than simply being put in by hand as an external source.

We now conclude the chapter on classical field theory by studying the simplest fundamental theory of such a kind.

Scalar Electrodynamics And Abelian Higgs Mechanism

The goal of this final subsection is to write down an interesting theory involving a complex scalar $\varphi$ and a vector field $A_\mu$. 'Interesting' is defined to mean that the two fields interact, which requires a Lagrangian involving terms with both $\varphi$ and $A_\mu$. The problem is now that in general $\varphi$ and $A_\mu$ do not transform trivially under a given symmetry. E.g., the complex scalar $\varphi$ transforms as (1.103) under a global $U(1)$, while $A_\mu$ obeys the gauge transformation (1.150).
The important point to notice is that both transformations involve only a single parameter, namely \( \alpha \) and \( f \). However, \( \alpha = \text{const.} \) is space-time independent, while \( f = f(x) \) is a function of the space-time coordinate \( x \). What happens if we generalize (1.103) by allowing the \( U(1) \) symmetry to be local rather than global? Our fields \( \varphi \) and \( \varphi^* \) then transform as

\[
\varphi \rightarrow \varphi' = e^{i\alpha(x)} \varphi, \quad \varphi^* \rightarrow (\varphi')^* = e^{-i\alpha(x)} \varphi^*.
\]

(1.175)

Since the scalar potential is polynomial in \( \varphi^* \varphi \) it is clearly invariant under these local \( U(1) \) transformations. But as already noted after (1.105), for the kinetic term to be invariant, we needed \( \alpha = \text{const.} \). Explicitly, we have

\[
\partial_\mu \varphi \rightarrow \partial_\mu \left(e^{i\alpha(x)} \varphi \right) = e^{i\alpha(x)} \left( \partial_\mu + i \hat{\partial}_\mu \alpha(x) \right) \varphi.
\]

(1.176)

This prompts us to realise that if we identify \( \alpha(x) = f(x) \) then

\[
D_\mu \varphi = (\partial_\mu - iA_\mu) \varphi,
\]

(1.178)

has the correct properties! It transforms as

\[
(D_\mu \varphi)^*(D^\mu \varphi) \rightarrow (D_\mu \varphi)^*(D^\mu \varphi),
\]

(1.179)

under

\[
\varphi \rightarrow e^{i\alpha} \varphi, \quad A_\mu \rightarrow A_\mu + \partial_\mu \alpha.
\]

(1.180)

In other words the new kinetic term \( (D_\mu \varphi)^*(D^\mu \varphi) \), with \( \partial_\mu \) replaced by the covariant derivative \( D_\mu \), is gauge invariant. Notice that the covariant derivative (1.178) has dimension \([D_\mu] = 1\) and involves a “minimal” coupling \( A_\mu \varphi \) between the vector field \( A_\mu \) and the scalar \( \varphi \).

We can now combine our scalar and gauge field Lagrangian to obtain

\[
\mathcal{L} = (D_\mu \varphi)^*(D^\mu \varphi) - V(\varphi^* \varphi) - \frac{1}{4} F_{\mu\nu} F^{\mu\nu}, \quad V(\varphi^* \varphi) = V_0 + m^2 \varphi^* \varphi + \frac{\lambda}{4} (\varphi^* \varphi)^2,
\]

(1.181)

which is the Lagrangian of scalar electrodynamics - describing a charged complex scalar coupled to electromagnetism. Notice that in this new theory the role of \( A_\mu \) is to facilitate the invariance of the scalar field theory under local gauge transformations. In fact, if we would have started from the global \( U(1) \) invariant scalar theory, with the task of finding its locally \( U(1) \) invariant version, we would have necessarily been led to introducing a gauge field \( A_\mu \).

We will now use (1.181) as a toy model to study the spontaneous symmetry breaking of a local \( U(1) \) symmetry. Since the scalar potential is identical to the one in the global case, we can reuse all our old results. The case \( m^2 \geq 0 \) is again not very interesting since the symmetry is unbroken, so let’s focus on the situation with \( m^2 < 0 \). In the latter case the potential takes
the form of a “Mexican hat” as in Figure 1.4. We parametrize $\varphi$ around the minima in the following way

$$\varphi = \frac{1}{\sqrt{2}} (v_0 + h) e^{ix},$$  \hspace{1cm} (1.182)

with $v_0$ given in (1.110). It is straightforward to see that $h$ corresponds to a real massive field, while $\chi$ is a Goldstone mode (also a real field). We can now eliminate the Goldstone mode from $\varphi$ by choosing the specific gauge

$$\varphi \rightarrow \varphi' = e^{-ix} \varphi = \frac{1}{\sqrt{2}} (v_0 + h).$$  \hspace{1cm} (1.183)

Under this transformation, we see from (1.180) that

$$A_\mu \rightarrow A'_\mu = A_\mu - \partial_\mu \chi.$$  \hspace{1cm} (1.184)

Since the Lagrangian (1.181) is gauge invariant, we can as well write it in terms of primed rather than unprimed fields. For the covariant derivative, one finds

$$(D_\mu \varphi)' = \frac{1}{\sqrt{2}} \left( \partial_\mu h - i (v_0 + h) A'_\mu \right),$$  \hspace{1cm} (1.185)

and inserting this expression into (1.181) leads to

$$\mathcal{L} = \frac{1}{2} (\partial_\mu h)^2 + \frac{1}{2} A'_\mu A'^\mu (v_0^2 + 2v_0 h + h^2) - V(h) - \frac{1}{4} F'_\mu F'^{\mu},$$  \hspace{1cm} (1.186)

with

$$V(h) = V_0 + \frac{1}{4} m^2 v_0^2 - m^2 h^2 + \frac{\lambda}{16} (4v_0 h^3 + h^4).$$  \hspace{1cm} (1.187)

The most striking feature of this result is that the Goldstone mode $\chi$ has (magically) disappeared and we are only left with the fields $h$ and $A'_\mu$. So did we lose a dof? No, because the vector field is now no longer massless, but has a mass

$$m_{A'} = v_0.$$  \hspace{1cm} (1.188)

It has hence three dofs (as opposed to just two in the massless case), which explains the disappearance of $\chi$: the Goldstone boson has been “eaten” by the vector field and provides the additional longitudinal component for the massive $A'_\mu$. This feature can also be inferred from (1.184). Notice furthermore that the mass of $h$ is given by

$$m_h = \sqrt{-2m^2},$$  \hspace{1cm} (1.189)

and that (1.186) and (1.187) have a rich structure of interactions: there are couplings of two gauge fields to one and two massive scalars ($A' A'$ and $A' A'^2$) as well as cubic and quartic self-couplings of the massive scalar ($h^3$ and $h^4$).

So let’s briefly summarize our main findings: we have learned that a spontaneously broken local (or gauge) symmetry leads to a mass for the associated vector boson and the conversion
of the Goldstone boson into the longitudinal component of the gauge field. This is the famous Higgs mechanism that we already alluded to.

The Higgs mechanism had already appeared within superconductivity. This can be understood as the spontaneous breaking of electromagnetism within a superconductor, so that the photon is effectively massive and the magnetic field is expelled from the interior of the superconductor.

The massive scalar $h$ left over in our model is called the Higgs scalar. The phenomenon described above (but in its generalization to non-abelian gauge groups) is at work in the SM of particle physics. In fact, since the SM symmetry breaking pattern is $SU(2)_L \times U(1)_Y \rightarrow U(1)_{em}$, there are three Goldstone bosons that give masses to three electroweak gauge bosons. The masses of these particles are all proportional to the symmetry breaking scale (or vacuum expectation value) $v_0$. A detailed discussion of the breaking of the electroweak symmetry, while worthwhile, is beyond the scope of this lecture, as it requires the introduction of non-abelian gauge symmetries, but you will learn a bit more in one of the exercises.