

Quantum Field Theory Notes

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NOTE DESCRIPTION

This note is author's self-study notes of quantum field theory and what I had chose is the *Quantum Field Theory and the Standard Model* written by Matthew D. Schwartz. In this note, I shall use the same names of the chapters to fit with the book itself and almostly there will exist more things about my thoughts which may have some mistakes.

pS: *And the version is the cloud version which run on overleaf because of the hard work of setting automatic English spelling in my local textstudio. So this can be regarded as my first attmpt on the cloud.*

*To those who taught me physics.
To human rationality.*

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Part I

Field Theory

Chapter 1

Microscopic Theory of Radiation

In this chapter, the Schwartz had shown to us that how the quantum theory was brought up especially the Plank's work in blackbody radiation also the Einstein's coefficients of spontaneous emission, and we shall use the first principle of quantum field theory to illustrate plenty of conclusions and phenomena. Also, we have to say, this is a period of history.

Let's go back to the classical time. We can derive the formula of frequencies from electrodynamics: the ω_n in a box of size L supported standing electromagnetic waves

$$\omega_n = \frac{2\pi}{L} |\vec{n}| c \quad (1.1)$$

By the classical equipartition theorem, we can come up the concept of blackbody, an object at fixed temperature whose internal structure we do not care about. And, how to say? Blackbodies should emit light equally in all modes with the intensity in volume of phase space:

$$I(\omega) = \frac{1}{V} \frac{d}{d\omega} E(\omega) = \text{constant} \times c^{-3} \omega^2 k_B T \quad (1.2)$$

But that will upset us, becasue the I is the quadratic function of ω . That says, when the ω grows up, the I will diverge. That's the Ultra-violet catastrophe.

Experimentally, the distrubution looks more like a Maxwell-Boltzmann distribution which imply that there is no diverge from the frequency interval $[0, \infty]$. In order to deal with the ultraviolet disaster, Planck successfully matched the experimental data with phenomenological methods such as interpolation, but it required the assumption that energy was quantized and introduced a new constant to distinguish the classical bewteen quantum mechanics:

$$E_n = \hbar \omega_n = \frac{2\pi}{L} \hbar |\vec{n}| = |\vec{p}_n| \quad (1.3)$$

Later, the Albert Einstein interpreted this as implying that light is made up of particles, the photons. Anyway, quantizing light resolves the blackbody paradox and light having energy leads to the photoelectric effect which was proven by Compton scattering experiment.

Anyway, with Planck's energy hypothesis, we can compute the specturm functions from first principle: each mode of frequency ω_n can be excited an integer number j times giving energy

$jE_n = j\hbar\omega$ in that mode. The partition function from Boltzmann distribution is $\exp(-\text{energy} \cdot \beta)$, thus the mean value of energy can be computed by followings:

$$\begin{aligned}\langle E_n \rangle &= \frac{\sum_{j=0}^{+\infty} jE_n e^{-jE_n\beta}}{\sum_{j=0}^{+\infty} e^{-jE_n\beta}} && \text{Convergent infinite series} \\ &= \frac{-\frac{d}{d\beta} \frac{1}{1-e^{-\hbar\omega_n\beta}}}{\frac{1}{1-e^{-\hbar\omega_n\beta}}} \\ &= \frac{\hbar\omega_n}{e^{\hbar\omega_n\beta} - 1}\end{aligned}\tag{1.4}$$

When we take the continue limit, namely $L \rightarrow \infty$, we turn the \sum to \int :

$$\begin{aligned}E(\omega) &= \int^\omega d^3\vec{n} \frac{\hbar\omega_n}{e^{\hbar\omega_n\beta} - 1} \quad \text{for all space} \\ &= 4\pi\hbar \frac{L^3}{8\pi^3} \int^\omega d\omega' \frac{(\omega')^3}{e^{\hbar\omega'\beta} - 1}\end{aligned}\tag{1.5}$$

so our first question have been figured out. Thus, we have

$$I(\omega) = \frac{1}{V} \frac{dE(\omega)}{d\omega} = \frac{\hbar}{\pi^2} \frac{\omega^3}{e^{\hbar\omega\beta} - 1}\tag{1.6}$$

So the UV catastrophe disappeared when we introduced quantization of energy!

What does this have to do with quantum field theory? In order for this derivation, which used equilibrium statistical mechanics, to make sence, light has to be able to equilibrate. For example, if we heat up a box with monochromatic light, eventually all frequencies must be excited. However, if different frequencies are different particles, equilibration must involve one kind of particle turning into another kind of particle. So, particles must be created and destroyed. Quantum field theory tells us how that happens!

That's exciting rational analysis!

However, there's one thing worth thinking. When we talk about the UV, we will find that it is a statistical physics problem. Yeah, we used to tackle the spontaneous emission through the statistical method. But the most exciting part is that we tackle the spontaneous emission through the single atom from first principle developed from quantum field theory not the statistical mechanics. In other words, Einstein's spontaneous emission coefficient can be calculated from the first principles of quantum field theory without the use of statistical mechanical conditions that assume thermal equilibrium. So this is a new evidence and defense to the quantum field theory. Let's see his face! (We're just going to talk about how do we get a solution to this problem from quantum field theory)

Before the start, we plain the results from Einstein's A,B coefficients:

$$\begin{aligned}B' &= B \\ \frac{A}{B} &= \frac{\hbar}{\pi^2} \omega^3\end{aligned}\tag{1.7}$$

The $B = B'$, says that the coefficient of absorption must be the same as the coefficient for stimulated emission. The coefficient B and B' can be computed in quantum mechanics using time-dependent perturbation theory with an external electromagnetic field. The ratio of A and B determines A . However, does spontaneous emission from an atom have anything to do with equilibrium of a gas since an atom radiates at the same rate no matter what is around it. That need the quantum field theory to figure out until 10 years after Einstein's calculation.

We know that, there exists some basic and profound motion mode, the $1/r$ and r^2 . The first implies the motion of celestial bodies and the second implies an oscillator which often represents the cyclic motion. From the quantum mechanics, we can say that the easiest way to develop a quantum harmonic oscillator is with creation and annihilation operators, a^\dagger and a , satisfying:

$$[a, a^\dagger] = 1 \quad (1.8)$$

and the number operator $\hat{N} = a^\dagger a$, which counts modes

$$\hat{N} |n\rangle = n |n\rangle \quad (1.9)$$

Notice that $a^\dagger |n\rangle$ is the eigenstate of the operator \hat{N} , proof are given:

$$\begin{aligned} \hat{N} a^\dagger |n\rangle &= a^\dagger a a^\dagger \quad \text{use the commutator } [a, a^\dagger] = a a^\dagger - a^\dagger a = 1 \\ &= a^\dagger (1 + a^\dagger a) |n\rangle \\ &= a^\dagger |n\rangle + a^\dagger a^\dagger a |n\rangle \\ &= a^\dagger |n\rangle + a^\dagger a^\dagger \sqrt{n} |n-1\rangle \\ &= a^\dagger |n\rangle + \sqrt{n} \sqrt{n} a^\dagger |n\rangle \\ &= (n+1) a^\dagger |n\rangle \end{aligned} \quad (1.10)$$

and the eigenvalue is $n+1$. Also, we normalize the eigenstates as $\langle n|n\rangle = 1$

Now, let's deal with the spontaneous emission. I think it is necessary to re-explain the concept of **spontaneous emission** quoted from Wikipedia.

Spontaneous emission is the process in which a quantum mechanical system (such as a molecule, an atom or a subatomic particle) transits from an excited energy state to a lower energy state (e.g., its ground state) and emits a quantized amount of energy in the form of a photon. Spontaneous emission is ultimately responsible for most of the light we see all around us; it is so ubiquitous that there are many names given to what is essentially the same process. If atoms (or molecules) are excited by some means other than heating, the spontaneous emission is called luminescence. For example, fireflies are luminescent. And there are different forms of luminescence depending on how excited atoms are produced (electroluminescence, chemiluminescence etc.). If the excitation is affected by the absorption of radiation the spontaneous emission is called fluorescence. Sometimes molecules have a metastable level and continue to fluoresce long after the exciting radiation is turned

off; this is called phosphorescence. Figurines that glow in the dark are phosphorescent. Lasers start via spontaneous emission, then during continuous operation work by stimulated emission.

As we developed in quantum mechanics, the Fermi's Golden Rules

$$\Gamma \sim |\mathcal{M}|^2 \delta(E_f - E_i) \quad (1.11)$$

implies the transition rate between two state where δ function enforces the conservation of energy (We will derive the quantum field theory version of cross-section and decay rate in chapter 5). The matrix element (tip: perturbation theory) is

$$\mathcal{M} = \langle f | H_{int} | i \rangle \quad (1.12)$$

namely the schematic

$$\mathcal{M} = \left\{ \text{final state} \xleftarrow[\text{Interaction}]{\text{Hamiltonian}} \text{initial state} \right\} \quad (1.13)$$

No matter what the H_{int} is exactly, it should have the creation and annihilation operator in line with the case of spontaneous emission and it must be Hermitian. Thus, we can write it formally without losing the physical meaning as following:

$$H_{int} = H_I^\dagger a^\dagger + H_I a \quad (1.14)$$

As for the process of spontaneous emission, we note the initial excited atom state as $|\text{Atom 2}; n_\omega\rangle$ with n_ω photons of frequency $\omega = \Delta/\hbar$:

$$|i\rangle = |\text{Atom 2}; n_\omega\rangle \quad (1.15)$$

and the final state is a lower energy atom we note as Atom 1 with $n_\omega + 1$ photons of energy Δ :

$$|f\rangle = |\text{Atom 1}; n_\omega + 1\rangle \quad (1.16)$$

and

$$\left\{ |\text{Atom 1}; n_\omega + 1\rangle \xleftarrow[\Delta=\hbar\omega]{\text{A Photon}} |\text{Atom 2}; n_\omega\rangle \right\} \quad (1.17)$$

Then we can compute:

$$\begin{aligned} \langle f | H_{int} | i \rangle &= \langle f | H_I^\dagger a^\dagger + H_I a | i \rangle \\ &= \langle \text{Atom 1}; n_\omega + 1 | H_I^\dagger a^\dagger + H_I a | \text{Atom 2}; n_\omega \rangle \\ &= \langle \text{Atom 1}; n_\omega + 1 | H_I^\dagger a^\dagger | \text{Atom 2}; n_\omega \rangle + \langle \text{Atom 1}; n_\omega + 1 | H_I a | \text{Atom 2}; n_\omega \rangle \\ &= \langle \text{Atom 1} | H_I^\dagger | \text{Atom 2} \rangle \langle n_\omega + 1 | a^\dagger | n_\omega \rangle + \langle \text{Atom 1} | H_I | \text{Atom 2} \rangle \langle n_\omega + 1 | a | n_\omega \rangle \\ &= \langle \text{Atom 1} | H_I^\dagger | \text{Atom 2} \rangle \sqrt{n_\omega + 1} \end{aligned} \quad (1.18)$$

which we have used the algebraic eigen-equation of operat a^\dagger , a and we decompose the eigenstates $|\text{Atom 1}; n_\omega + 1\rangle$, $|\text{Atom 2}; n_\omega\rangle$ and compound operators $H_I^\dagger a^\dagger$, $H_I a$ here. Further, we note that $\mathcal{M}_0^\dagger = \langle \text{Atom 1} | H_I^\dagger | \text{Atom 2} \rangle$. Thus, the process can be described as

$$|\mathcal{M}_{2 \rightarrow 1}|^2 = |\mathcal{M}_0|^2 (n_\omega + 1) \quad (1.19)$$

Instead, we excite an atom, then the initial state is unexcited atom with n_ω photons. Then, the initial and final state can be written as

$$|i\rangle = |\text{Atom 1}; n_\omega\rangle \quad (1.20)$$

$$|f\rangle = |\text{Atom 2}; n_\omega - 1\rangle \quad (1.21)$$

Same as (1.18), we have:

$$\begin{aligned} \langle f | H_{int} | i \rangle &= \langle \text{Atom 2} | H_I^\dagger | \text{Atom 1} \rangle \langle n_\omega - 1 | a^\dagger | n_\omega \rangle + \langle \text{Atom 2} | H_I | \text{Atom 1} \rangle \langle n_\omega - 1 | a | n_\omega \rangle \\ &= \langle \text{Atom 2} | H_I | \text{Atom 1} \rangle \sqrt{n_\omega} \end{aligned} \quad (1.22)$$

and

$$\mathcal{M}_{1 \rightarrow 2} = \mathcal{M}_0 \sqrt{n_\omega} \quad (1.23)$$

namely

$$|\mathcal{M}_{1 \rightarrow 2}|^2 = |\mathcal{M}_0|^2 n_\omega \quad (1.24)$$

According to Einstein's theory, the density of number satisfies (see on [Wikipedia](#))

$$\begin{aligned} dn_2 &= -dn_1 = -|\mathcal{M}_{2 \rightarrow 1}|^2 n_2 + |\mathcal{M}_{1 \rightarrow 2}|^2 n_1 \\ &= -|\mathcal{M}_0|^2 (n_\omega + 1) n_2 + |\mathcal{M}_0|^2 n_\omega n_1 \end{aligned} \quad (1.25)$$

while the original form are

$$dn_2 = -dn_1 = -[A + BI(\omega)]n_2 + B'I(\omega)n_1 \quad (1.26)$$

and we need to relate the number of photon modes of frequency ω to the intensity $I(\omega)$ since the energies are quantized by $\Delta = \hbar\omega = \hbar 2\pi/L|\vec{n}|$. The total energy is

$$E(\omega) = \int^\omega d^3\vec{n} \hbar\omega n_\omega = 4\pi\hbar L^3 \int^\omega \frac{d\omega}{(2\pi)^3} \omega^3 n_\omega \quad (1.27)$$

But we should multiply this by 2 for the two polarizations of light. Including the 2 factor, we can get the intensity

$$I(\omega) = \frac{1}{L^3} \frac{dE(\omega)}{d\omega} = \frac{\hbar\omega^3}{\pi^2} n_\omega \quad (1.28)$$

Thus, we have the exciting results:

$$dn_2 = -dn_1 = -|\mathcal{M}_0|^2 \left[1 + \frac{\pi^2}{\hbar\omega^3} I(\omega)\right] n_2 + |\mathcal{M}_0|^2 \left[\frac{\pi^2}{\hbar\omega^3} I(\omega)\right] n_1 \quad (1.29)$$

where we can read

$$B' = B, \quad \frac{A}{B} = \frac{\hbar}{\pi^2} \omega^3 \quad (1.30)$$

without any assumption of thermal equilibrium! That's one of the successes of Quantum Field Theory!

In a word, that's a history or beginning of our quantum field theory journey. More details can be found on my own paper notebook which follows the online course of Yu Jia from UCAS on [bilibili](#) starting from the Dirac and K-G fields.

Chapter 2

Lorentz Invariance and the Second Quantization

In the last chapter, we already know that treating each mode of electromagnetic radiation in a cavity as a simple harmonic oscillator, we can derive Einstein's relation between the coefficients on induced and spontaneous emission without resorting to statistical mechanics. In this chapter, we gonna to talk about the scnd quantization and free field used to describe the canonical quantization of relativistic fields.

In M.D. Schwartz's book, however, he devotes a larger part of the book to a review of special relativity and some tensor operations, a part I will omit in this note. In the course on general relativity (Canbin Liang's textbook), we did a lot of tensor operations and discussion of the differential geometry of spacetime geometry, so we will not go into some metrics, i.e. the so-called covariance and inversion, time-like and spacelike concepts, etc., and will just use them directly instead of elaborating on them again in this note.

Here, it is important to recall that there are two approaches to developing quantum field theory, canonical quantization and Feynman integrals. Here, we will mainly focus on the analogy of harmonic oscillators then slowly build up the concept of canonical quantization. As for the Feynman integral, we will briefly introduce its concepts as well as operations, which mainly refer to the books of A.Zee ,*Quantum field theory in a nutshell*¹(whose said that from our point of view, quantum mechanics is nothing more than 0+1 dimensional quantum field theory).

to be supplemented

¹Very great book!

Chapter 3

Classical Field Theory

As we had illustrated formerly that the quantum field theory is just quantum mechanics with as infinite number of oscillators. We already saw that it can do some remarkable things, such as explain spontaneous emission without the statistical theory. (But it also seems to lead to absurdities, such as an infinite shift in the energy levels of hydrogen atom). Remember that quantum field theory is not absurd but extremely predictive. That means you should be very careful about how we do calculations that the theory gets right without infinities. These are called the *tree-level processes*, which means they are leading order in an expansion in \hbar . Since taking $\hbar \rightarrow 0$ gives the classical limit, tree-level calculations are closely related to calculations in classical field theory, which is the subject of this chapter.

A classical field theory is just a mechanical system with a continuous set of degrees of freedom. Differently, we often use the density of Hamiltonians or Lagrangians in place of Hamiltonian or Lagrangians: (notice that we often omit the word *density*, so that whenever we refer to Hamiltonians and Lagrangians we refer to their density)

$$H = \int d^3x \mathcal{H}, \quad L = \int d^3x \mathcal{L} \quad (3.1)$$

and the dynamics for a Lagrangian system are determined by *the principle of least action*. The action is the integral over the time:

$$S = \int dt L = \int d^4x \mathcal{L} \quad (3.2)$$

Formally, the Hamiltonian is a functional of fields and their conjugate momenta $\mathcal{H} = \mathcal{H}[\phi, \pi]$ while the Lagrangian is the *Legendre transform of the Hamiltonian* defined as

$$\mathcal{L}[\phi, \dot{\phi}] = \pi[\phi, \dot{\phi}] \dot{\phi} - \mathcal{H}[\phi, \pi[\phi, \dot{\phi}]] \quad (3.3)$$

where items satisfies

$$\dot{\phi} = \partial_t \phi$$

and $\pi[\phi, \dot{\phi}]$ is defined by

$$\frac{\partial \mathcal{H}[\phi, \pi]}{\partial \pi} = \dot{\phi}$$

The inverse transform is

$$\mathcal{H}[\phi, \pi] = \pi \dot{\phi}[\phi, \pi] - \mathcal{L}[\phi, \dot{\phi}[\phi, \pi]] \quad (3.4)$$

where the item $\dot{\phi}[\phi, \pi]$ is defined by

$$\frac{\partial \mathcal{L}[\phi, \dot{\phi}]}{\partial \dot{\phi}} = \pi$$

What a coup is that you can regard the ϕ as the analog of q and regard the π as the analog of p which represents that the thought of field is the field version of the so-called analytical mechanics.

Also, we can derive the Hamiltonian with the sum of the kinetic and potential energies of a system:

$$\mathcal{H} = \mathcal{K} + \mathcal{V} \quad (3.5)$$

while the Lagrangian was shown as:

$$\mathcal{L} = \mathcal{K} - \mathcal{V} \quad (3.6)$$

The Hamiltonian corresponds to a conserved quantity that is the total energy of a system while the Lagrangian does not. The question is that, Hamiltonians, however, is that they are not Lorentz invariant. Although the Hamiltonian picks out the energy or precisely the energy eigenvalue, which is not a Lorentz scalar. On the contrary, it is the 0 component of a Lorentz vector $P_\mu = (H, \vec{p})$ (4-momentum vector). The Hamiltonian density is the 00 component of a Lorentz tensor, the energy-momentum tensor $\mathcal{T}_{\mu\nu}$ which we will see it is the Noether current given by symmetry under global space-time translations.

What's more? The Matthew D. Schwartz taught us a fact or a convention. Hamiltonians are great for non-relativistic systems (such as we have already learnt in quantum mechanics), but for relativistic systems we will almost exclusively use Lagrangians.

And there are many terms we should understand clearly. We do not usually talk about kinetic and potential energy in quantum field theory. Instead we talk about *kinetic terms* and then about *interactions*, for reasons that will become clear after we have done a few calculations (shut up and do calculations!). *Kinetic terms* are *bilinear*, meaning they have exactly two fields as we saw the square item in Newton's mechanics. So kinetic terms may look like:

$$\mathcal{L}_K \supset \frac{1}{2} \phi \square \phi, \quad \bar{\psi} \not{\partial} \psi, \quad \frac{1}{4} F_{\mu\nu}^2, \quad \frac{1}{2} m^2 \phi^2, \quad \frac{1}{2} \phi_1 \square \phi_2$$

Anything with just two fields of the same or different type can be called a kinetic term. The kinetic terms tell you about the free (non-interacting) behavior. Fields with kinetic terms are said to be *dynamical* or *propagating*. More precisely, a field should have time derivatives in its kinetic term to be dynamical. It is also sometimes useful to think of a *mass term*, such as $m^2 \phi^2$, as an interaction rather than a kinetic term. *Interactions* have three or more fields like $\mathcal{L} \supset \lambda \phi^3, g \bar{\psi} A \psi, g \partial_\mu \phi A_\mu \phi^*, g^2 A_\mu^2 A_\nu^2, 1/M_{PI} \partial_\mu h_{\mu\nu} \partial_\nu h_{\alpha\beta} h_{\alpha\beta}, \dots$ Since the interactions are everything but the kinetic terms, we also sometimes write $\mathcal{L} = -\mathcal{V} = -\mathcal{H}_{int}$. It is helpful if the coefficients of the interaction terms are small in some sense, so that the fields are *weakly interacting* and we can do *perturbation theory*.

From now, we focus on the principal of least action and consider the varying of $\phi \rightarrow \phi + \delta\phi$, and the Lagrangian is the functional of a field and its first derivatives¹, then we can do functional computation as following:

$$\begin{aligned}
\delta S &= \int d^4x [\mathcal{L} + \delta\mathcal{L} - \mathcal{L}] \\
&= \int d^4x \left[\frac{\partial\mathcal{L}}{\partial\phi} \delta\phi + \frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \delta(\partial_\mu\phi) \right] \\
&= \int d^4x \left[\frac{\partial\mathcal{L}}{\partial\phi} \delta\phi + \frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \partial_\mu(\delta\phi) \right] \\
&= \int d^4x \left[\frac{\partial\mathcal{L}}{\partial\phi} \delta\phi + \partial_\mu \left(\frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \delta\phi \right) - \left(\partial_\mu \frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \right) \delta\phi \right] \\
&= \int d^4x \left\{ \left[\frac{\partial\mathcal{L}}{\partial\phi} - \partial_\mu \frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \right] \delta\phi + \partial_\mu \left[\frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \delta\phi \right] \right\} \\
&= \left\{ \int d^4x \left[\frac{\partial\mathcal{L}}{\partial\phi} - \partial_\mu \frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \right] \delta\phi \right\} + \left[\frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \delta\phi \right] \Big|_{\text{boundaries}} \\
&= \int d^4x \left[\frac{\partial\mathcal{L}}{\partial\phi} - \partial_\mu \frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} \right] \delta\phi
\end{aligned} \tag{3.7}$$

Above the computation is the variation of the action², and we should note that the principle of the least action is beyond the classical mechanics as the equations of motion determined by this principle in classical field theory. Namely it should be insensitive to small variations of those field $\delta S/\delta\phi = 0$, and if this holds for all variations, it gives $\delta S = 0$, thus:

$$\boxed{\frac{\partial\mathcal{L}}{\partial\phi} - \partial_\mu \frac{\partial\mathcal{L}}{\partial(\partial_\mu\phi)} = 0} \tag{3.8}$$

These are the *Euler-Lagrange equations* while they give the *equations of motion* following from a Lagrangian.

Let's check the principle from two examples. Given an action as the following:

$$S = \int d^4x \mathcal{L} = \int d^4x \left[\frac{1}{2} (\partial_\mu\phi)(\partial_\mu\phi) - \mathcal{V}[\phi] \right]$$

and we put the Lagrangian into the E-L equations and we can easily get the motion equation:

$$-\mathcal{V}'[\phi] - \partial_\mu(\partial_\mu\phi) = 0$$

where we often note the $\square = \partial_\mu^2$ as the *d'Alembertian*.

What an important one is the Lagrangian as the form as

$$\mathcal{L} = \frac{1}{2} (\partial_\mu\phi)(\partial_\mu\phi) - \frac{1}{2} m\phi^2 \tag{3.9}$$

¹The form $\mathcal{L}[\phi, \partial_\mu\phi]$ is the form that "classical" Lagrangians had. If only first derivatives are involved, boundary conditions can be specified by initial positions and velocities only in accordance with Newton's idea. And actually, the complex coupling fields items in kinetic or interactions must occur due to the quantum effects in all but the simplest *renormalizable* field theories such as generic in *effective field theories*.

²While computing, we use the method of integration by parts and we exchange the variation and derivative which we often do in physics.

When we put it in the E-L equations, it will give an important equation as :

$$\boxed{(\square + m^2)\phi = 0} \quad (3.10)$$

The equation (3.10) is known as the *Klein-Gordon equation* which describes the equation of motion for a *free scalar field*.

When we meet a physical quantity, a natural question is that will the quantity give some symmetry or what kind of symmetry is inside? Now, we will show the topic of syemmetry *Noether's theorem*.

We must clarify the concept that the symmetries we are referring to are not the "folded" symmetries obtained by folding a piece of paper , like the symmetries learned by schoolchildren³. The so-called symmetry in physics is actually a kind of invariance, under a certain transformation, the physical quantity we study does not change, then it is said that this is a certain symmetry.

As for our physical quantity, the Lagranian, it may be invariant under some special type of variation $\phi \rightarrow \phi + \delta\phi$. For example, given a Lagrangian for a complex field ϕ as:

$$\mathcal{L} = |\partial_\mu \phi|^2 - m^2 |\phi|^2 \quad (3.11)$$

it is the square, so when we plug the complex phrase in the state ϕ it will be invariant. Under transformation $\phi \rightarrow e^{-i\alpha} \phi$ for any $\alpha \in \mathbb{R}$, the Lagrangian is invariant. We claim that this transformation is a *symmetry* of the Lagrangian. Notice that the Lagrangian has two independent real degrees of freedom in such a complex field ϕ , actually you can take the $\phi = \phi_1 + i\phi_2$ then the Lagrangian is determined by two fields ϕ_1, ϕ_2 .

To fit our symmetry, the parameter α is continuous while we can deal the degrees of freedom as the kind of special discrete symmetry, namely the conjugate ϕ and ϕ^* . Consider that our example is as the form as the scalar free field which should fit the Klein-Gordon equation mentioned above.

We use a *continuous* parameter to describe the transformation and the symmetry which exactly mean the Lagrangian \mathcal{L} is invariant under the transformation by parameter α saying that the functional derivative of \mathcal{L} with respect to parameter α is 0, thus:

$$\frac{\delta \mathcal{L}}{\delta \alpha} = 0 \quad (3.12)$$

We have already compute the δL in the (3.7), then we can quickly write the $\delta \mathcal{L}/\delta \alpha$:

$$\frac{\delta \mathcal{L}}{\delta \alpha} = \sum_n \left\{ \left[\frac{\partial \mathcal{L}}{\partial \phi_n} - \partial_\mu \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_n)} \right] \frac{\delta \phi_n}{\delta \alpha} + \partial_\mu \left[\frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_n)} \frac{\delta \phi_n}{\delta \alpha} \right] \right\} = 0 \quad (3.13)$$

where ϕ_n are the set of fields the Lagrangian depends on⁴.

³Although folding a piece of paper does have a spatial symmetry

⁴Or you can regard ϕ_n as the ϕ and ϕ^* in our example mentioned above

When the equations of motion are satisfied⁵, the (3.13) will be written as:

$$\sum_n \partial_\mu \left[\frac{\partial \mathcal{L}}{\partial(\partial_\mu \phi_n)} \frac{\delta \phi_n}{\delta \alpha} \right] = 0 \quad (3.14)$$

implies the important *conservative equation*:

$$\boxed{\partial_\mu J_\mu = 0} \quad (3.15)$$

where

$$\boxed{J_\mu = \sum_n \frac{\partial \mathcal{L}}{\partial(\partial_\mu \phi_n)} \frac{\delta \phi_n}{\delta \alpha}} \quad (3.16)$$

is known as a *Noether Current*.

The Noether current can be used to describe the symmetry. For example, we can compute the Noether current from the example Lagrangian above $\mathcal{L} = |\partial_\mu \phi|^2 - m^2 |\phi|^2$. We note that $\phi_1 = \phi$, $\phi_2 = \phi^*$, and their variant derivative $\delta \phi / \delta \alpha = -i\phi$, $\delta \phi^* / \delta \alpha = i\phi^*$, so that

$$J_\mu = i(\phi^* \partial_\mu \phi - \phi \partial_\mu \phi^*) \quad (3.17)$$

and note that the symmetry is continuous so that we can take small variations. As this is the scalar free field, we can check that

$$\partial_\mu J_\mu = i(\phi^* \square \phi - \phi \square \phi^*) = im^2(\phi \phi^* - \phi^* \phi) = 0 \quad (3.18)$$

which vanishes when they are satisfied with the equations of motion⁶.

In the previous paragraph, we should note that the equation of motion is satisfied, and in this case vector field J_μ that satisfies $\partial_\mu J_\mu = 0$ is called a *conserved current*. It is called *conserved* because it is actually the *charge continuity equation*. Here, the total charge Q , defined as

$$Q = \int d^3x J_0 \quad (3.19)$$

satisfies

$$\partial_t Q = \int d^3x \partial_t J_0 = \int d^3x \nabla \cdot \vec{J} = 0 \quad (3.20)$$

The equation above is similar to the "classical" charge conserved, like the equation $\partial \rho / \partial t + \nabla \cdot (\rho \vec{v}) = 0$ in the fluid mechanics⁷ which called the Euler equation⁸. Also, we have assumed \vec{J} vanishes at the spatial boundary which is reasonable in physics, since, by assumption, nothing is leaving our experiment. That says, the total charge does not change with time, and is conserved.

So now, we can formally explain what *Noether's theorem* is:

⁵In contrast to (3.7), the (3.13) holds all ϕ_n because it is derived by continuous symmetry's computation of functional derivative so the field configurations ϕ_n can not fit in the equation of motion corresponds to a symmetry

⁶We used the equations $\square \phi = -m^2 \phi$ and $\square \phi^* = -m^2 \phi^*$.

⁷To fit with the total charge, you can rewrite the equation through the Gauss's Law

⁸You may notice that the $\partial_\mu J_\mu$ shows that $\partial_t J_0 - \nabla \cdot \vec{J}$ has a minus sign difference. It doesn't matter because of the difference of definition from the J_0 or the total charge Q

If a Lagrangian has a continuous symmetry then there exists a current associated with that symmetry that is conserved when the equations of motion are satisfied.

Notice that the symmetry must be continuous otherwise we cannot induce the Noether's theorem because the small variations $\delta\mathcal{L}/\delta\alpha$ has no meaning which also can be expressed that Noether's theorem does not apply to discrete symmetries⁹. And now we can summarize a few points¹⁰:

1. The symmetry must be continuous, otherwise $\delta\alpha$ has no meaning.
2. The current is conserved on-shell, that is, when the equations of motion are satisfied.
3. It works for global symmetries, parametrized by numbers α , not only for local (gauge) symmetries parametrized by functions $\alpha(x)$.

We have expounded the symmetry especially the Noether's theorem. We deal with the Noether current J_μ . But now we deal with an important case which is the global symmetry to the action not the Lagrangian¹¹. That's the global space-time translations and the symmetry gives a conserved current, the *energy-momentum tensor* $\mathcal{T}_{\mu\nu}$.

We take a translation from the whole space-time manifold, and we say that physics at the point x should be same as the point y ¹². We note the translation as an infinitesimal 4-vector ξ^μ , then for the scalar field, we have:

$$\phi(x) \rightarrow \phi(x + \xi) \quad (3.21)$$

Using the field expansion, we can expand the equation as

$$\phi(x + \xi) = \phi(x) + \sum_{\nu} \frac{\partial\phi(x)}{\partial x^\nu} \xi^\nu + \dots = \phi(x) + \xi^\nu \partial_\nu \phi(x) + \dots \quad (3.22)$$

That implies that the infinitesimal translation means:

$$\frac{\delta\phi}{\delta\xi^\nu} = \partial_\nu \phi \quad (3.23)$$

⁹Such as the symmetry under $\phi \rightarrow -\phi$ of $\mathcal{L} = 1/2\phi\Box\phi - m^2\phi^2 - \lambda\phi^4$ with ϕ real even though it is exactly variant under the transformation.

¹⁰There is no Noether's theorem in general relativity, but we can derive the global conservation in GR instead of the local conservation.

¹¹Notice that Noether's theorem fits to the continuous symmetry to the Lagrangians, but we deal with the case including the time part's translations so this is the kind of symmetry to action, naturally the Noether's theorem is not working.

¹²Notice that this is a translation of spacetime as a whole and not a change in the coordinate system.

namely as the Lagrangian:

$$\frac{\delta \mathcal{L}}{\delta \xi^\nu} = \partial_\nu \mathcal{L} \quad (3.24)$$

Since this is a total derivative, then the variation of action S :

$$\begin{aligned} \delta S &= \int d^4x \delta \mathcal{L} \\ &= \int d^4x \xi^\nu \partial_\nu \mathcal{L} = 0 \end{aligned} \quad (3.25)$$

which is why we sometimes say this is a symmetry of the action, not the Lagrangian.

And what's the conserved current? As we have computed in (3.13), and use the equations of motion, then the variation of the Lagrangian form a translation ξ^ν can be written as:

$$\frac{\partial \mathcal{L}}{\partial \xi^\nu} = \partial_\mu \left(\sum_n \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_n)} \frac{\delta \phi_n}{\delta \xi^\nu} \right) \quad (3.26)$$

using the (3.23), we have an equation:

$$\partial_\mu \left(\sum_n \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_n)} \partial_\nu \phi_n \right) = \partial_\nu \mathcal{L} \quad (3.27)$$

namely:

$$\partial_\mu \left(\sum_n \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_n)} \partial_\nu \phi_n - g_{\mu\nu} \mathcal{L} \right) = 0 \quad (3.28)$$

which is the conserved equation $\partial_\mu \mathcal{T}_{\mu\nu} = 0$ producing four Noether currents, and the $\mathcal{T}_{\mu\nu}$ is:

$$\boxed{\mathcal{T}_{\mu\nu} = \sum_n \frac{\partial \mathcal{L}}{\partial (\partial_\mu \phi_n)} \partial_\nu \phi_n - g_{\mu\nu} \mathcal{L}} \quad (3.29)$$

The four Noether currents correspond to four conserved quantities, the energy and the momentum. $\mathcal{T}_{\mu\nu}$ is called the *energy-momentum tensor*. We see the 00 component of the energy-momentum tensor:

$$\varepsilon = \mathcal{T}_{00} = \sum_n \frac{\partial \mathcal{L}}{\partial \dot{\phi}_n} \dot{\phi}_n - \mathcal{L} \quad (3.30)$$

which has the same form as the Legendre transform of Lagrangian for the $\partial \mathcal{L} / \partial \dot{\phi}_n = \pi_n$. Then the 00 component of energy-momentum tensor is the energy density. By the way, the definition of energy-momentum tensor in (3.29) is called the *canonical energy-momentum tensor* for the difference definition from the general relativity¹³.

There's an another thing we should clarify. The concept of the current is abstract current which we use it to describe the continuous symmetry. And with that comes the concept of the abstract charge. Both the conserved vector J_μ associated with a global symmetry and the energy-momentum tensor $\mathcal{L}_{\mu\nu}$ are types of currents. The currents are used in many ways in quantum field theory, for example:

¹³There is another way to derive the energy-momentum tensor in general relativity and Lagrangian form in GR. See textbook *Introduction to differential geometry and general relativity* written by Canbin Liang and Bin Zhou

-
1. Currents can be Noether currents associated with a symmetry.
 2. Currents can refer to *external* currents¹⁴.
 3. Currents can be used as sources for fields appearing in the Lagrangians.
 4. Currents can be place-holders for certain terms in a Lagrangian.

Now, we gonna to do some calculations, the Coulomb's law and Coulomb potential, using the classical field theory we developed above, and we will induce the method of Green's function and the Feynman diagrams (rules) in these examples. We start with an external current:

$$J_\mu(x) = \begin{cases} J_0(x) = \rho(x) = e\delta^3(x) \\ J_i(x) = 0 \end{cases} \quad (3.31)$$

whose Lagrangian¹⁵ is:

$$\mathcal{L} = -\frac{1}{4}F_{\mu\nu}^2 - A_\mu J_\mu \quad (3.32)$$

and the E-L equations¹⁶ imply the :

$$\partial_\mu F_{\mu\nu} = J_\nu \quad (3.33)$$

which is the Maxwell equations and choose the Lorenz gauge, $\partial_\mu A_\mu = 0$, then:

$$\square A_\nu(x) = J_\nu(x) \quad (3.34)$$

and the formal solution :

$$A_\nu(x) = \frac{1}{\square} J_\nu(x) \quad (3.35)$$

The $1/\square$ just means the inverse of the \square . It says that the A_ν field is determined by the source J_ν after it propagates with the *propagator*. It's an insightful way to understand.

$$\Pi_A = \frac{1}{\square} \quad (3.36)$$

Expect the electromagnetic waves' solutions, the point charge at the origin will gives:

$$\begin{aligned} A_i &= 0 \\ A_0(x) &= \frac{e}{\square} \delta^3(x) \end{aligned} \quad (3.37)$$

To solve the equations, we develop the tools of fourier transform first. We have the general transform:

$$\delta^3(\vec{x}) = \int \frac{d^3k}{(2\pi)^2} e^{i\vec{k}\vec{x}} \quad (3.38)$$

¹⁴These are given background configurations, such as electrons flowing through a wire

¹⁵Please look up the textbook of electrodynamics and write the Lagrangian of electromagnetic field

¹⁶Using the symmetry of $F_{\mu\nu}$

Since the Laplacian is $\Delta = \partial^2$ for space parts, therefore:

$$\begin{aligned}
\Delta \delta^3(\vec{x}) &= \Delta \int \frac{d^3k}{(2\pi)^3} e^{i\vec{k}\vec{x}} \\
&= (\partial_x^2 + \partial_y^2 + \partial_z^2) \int \frac{d^3k}{(2\pi)^3} e^{ik_x x} e^{ik_y y} e^{ik_z z} \\
&= i^2(k_x^2 + k_y^2 + k_z^2) \int \frac{d^3k}{(2\pi)^3} e^{i\vec{k}\vec{x}} \\
&= (i\vec{k})^2 \int \frac{d^3k}{(2\pi)^3} e^{i\vec{k}\vec{x}}
\end{aligned}$$

which we can extend the relations as following:

$$\begin{aligned}
\Delta^2 \delta^3(\vec{x}) &= (i\vec{k})^{2n} \int \frac{d^3k}{(2\pi)^3} e^{i\vec{k}\vec{x}} \\
&= -(\vec{k})^2 \int \frac{d^3k}{(2\pi)^3} e^{i\vec{k}\vec{x}}
\end{aligned} \tag{3.39}$$

or we can rewrite the equation as the fourier transform equation, that means:

$$[\Delta^2 \delta^3(\vec{k})]^{-1} = -(\vec{k})^2 \tag{3.40}$$

Now, let's check the d'Alembertian for a lorentz-invariant quantities:

$$\delta^4(x) = \int \frac{d^4k}{(2\pi)} e^{ikx} \tag{3.41}$$

where the item kx means the summarization $k_0 x_0 - \vec{k} \cdot \vec{x}$. Therefore,

$$\square \delta^4(x) = - \int \frac{d^4k}{(2\pi)} e^{ikx} k^2 e^{ikx} \tag{3.42}$$

thus, the extension:

$$\square^n \delta^4(x) = \int \frac{d^4k}{(2\pi)^4} (-k^2)^n e^{ikx} \tag{3.43}$$

The fourier's law tells us that every functions can be expanded by fourier transform:

$$f(x) = \int \frac{d^4k}{(2\pi)^4} \tilde{f}(k) e^{ikx} \tag{3.44}$$

then

$$\square^n f(x) = \int \frac{d^4k}{(2\pi)^4} \square^n \tilde{f} e^{ikx} = \int \frac{d^4k}{(2\pi)^4} (-k^2)^n \tilde{f}(k) e^{ikx} \tag{3.45}$$

or namely

$$\widetilde{[\square^n f]}(k) = (-k^2)^n \tilde{f}(k) \tag{3.46}$$

.

Then we develop the important tools for Δ and \square which for a field theorist, they mean¹⁷:

$$\Delta \rightleftharpoons -\vec{k}^2, \quad \square \rightleftharpoons -k^2 \quad (3.47)$$

are the linear and it will simplify our computation.

So, let's back to our example, the equation:

$$A_0(x) = \frac{e}{\square} \delta^3(\vec{x}) = -\frac{e}{\Delta} \delta^3(\vec{x}) \quad (3.48)$$

use our tools, then we will have the integral in k -space¹⁸:

$$\begin{aligned} A_0(x) &= -\frac{e}{\Delta} \delta^3(\vec{x}) \\ &= \frac{e}{(2\pi)^3} \int \frac{d^3k}{\vec{k}^2} e^{i\vec{k}\vec{x}} \\ &= \frac{e}{(2\pi)^3} \int d^3k \frac{e^{i\vec{k}\vec{x}}}{k^2} \\ &= \frac{e}{(2\pi)^3} \iiint k^2 \sin\theta d\theta d\varphi dk \frac{e^{ikr \cos\theta}}{k^2} \\ &= \frac{e}{(2\pi)^3} \int_0^\infty k^2 dk \int_0^\pi \sin\theta d\theta \int_0^{2\pi} d\varphi \frac{e^{ikr \cos\theta}}{k^2} \\ &= \frac{e}{8\pi^2} \frac{1}{ir} \int_{-\infty}^{+\infty} dk \frac{e^{ikr} - e^{-ikr}}{k} \end{aligned} \quad (3.49)$$

and we use the method of mathematics in physics, the theory of complex variables functions,

$$\int_{-\infty}^{+\infty} dk \frac{e^{ikr} - e^{-ikr}}{k} = \lim_{\delta \rightarrow 0} \left[\int_{-\infty}^{+\infty} dk \frac{e^{ikr} - e^{-ikr}}{k + i\delta} \right] = \lim_{\delta \rightarrow 0} 2\pi i e^{-\delta r} = 2\pi i \quad (3.50)$$

therefore the equation

$$A_0(x) = \frac{e}{4\pi r} \quad (3.51)$$

shows the Coulomb potential.

From the important example above, we can summarize a kind of equations whose form follow the form:

$$\text{Field} = \frac{1}{\square} \text{Current} \quad (3.52)$$

We add some items in the Lagrangian \mathcal{L} . We consider the Lagrangian of a charged object radiating the A field,

$$\mathcal{L} = -\frac{1}{4} F_{\mu\nu}^2 - \phi^* \square \phi - ie A_\mu (\phi^* \partial_\mu \phi - \phi \partial_\mu \phi^*) \quad (3.53)$$

We can get the current in the space-time is $J_\mu = ie(\phi^* \partial_\mu \phi - \phi \partial_\mu \phi^*)$ and we can get the same equation in (3.52) in Lorenz gauge.

¹⁷In electrodynamics, we used to do the same with derivative operators, for example, by treating the derivative operator directly as $\nabla = i\vec{k}$ and the time derivative operator as $\partial_t = -i\omega$ for the plane wave solution or low order spherical wave solution.

¹⁸use the correct integral on the spherical surface and use the general integrand. Just one step by step

Using the propagators is a very useful way to solve these types of equations, and quite general. And we gonna to introduce an example about the graviton and we shall use the example's perturbation solution to induce the Green's function method.

Given the Lagrangian:

$$\mathcal{L} = -\frac{1}{2}h\Box h + \frac{1}{3}\lambda h^3 + Jh \quad (3.54)$$

where the h represents the gravitational potential. As the E-L equation (3.8), we get the equation of motion:

$$\Box h - \lambda h^2 - J = 0 \quad (3.55)$$

Solve the perturbatively in λ , for the $\lambda = 0$, the zero solution is

$$h_0 = \frac{1}{\Box}J \quad (3.56)$$

. Then we plug in

$$h = h_0 + h_1 \quad (3.57)$$

with $h_1 = \mathcal{O}(\lambda^1)$. Then

$$\Box(h_0 + h_1) - \lambda(h_0 + h_1)^2 - J = 0 \quad (3.58)$$

implies

$$\Box h_1 = \lambda h_0^2 + \mathcal{O}(\lambda^2) \quad (3.59)$$

so that:

$$h_1 = \lambda \frac{1}{\Box}(h_0 h_0) = \lambda \frac{1}{\Box}\left[\left(\frac{1}{\Box}J\right)\left(\frac{1}{\Box}J\right)\right] \quad (3.60)$$

So the solution to order λ in perturbation theory is

$$h = \frac{1}{\Box}J + \lambda \frac{1}{\Box}\left[\left(\frac{1}{\Box}J\right)\left(\frac{1}{\Box}J\right)\right] + \mathcal{O}(\lambda^2) \quad (3.61)$$

which can keep iterating.

This is the *Green's function method*, and the

$$\Pi = \frac{1}{\Box} \quad (3.62)$$

is called *2-point Green's function* or *propagator*. As its name illustrates, the propagator is a kind of tools to describe the propagating having nothing to do with the source exactly determined by the kinetic terms for a field which is insightful to physical pictures.

About this expansion, we can define the propagator as the solution to¹⁹ :

$$\Box_x \Pi(x, y) = -\delta^4(x - y) \quad (3.63)$$

Drop out the boundaries in infinite space, we can get

$$\Pi(x, y) = \int \frac{d^4 k}{(2\pi)^4} e^{ik(x-y)} \frac{1}{k^2} \quad (3.64)$$

¹⁹See the textbook *Methods of mathematics in physics*

The 2-point propagator is symmetry,

$$\Pi(x, y) = \Pi(y, x) \quad (3.65)$$

Using the $\square_y \Pi(x, y) = -\delta^4(x, y)$, we can write a field as

$$\begin{aligned} h(x) &= \int d^4y \delta^4(x - y) h(y) \\ &= - \int d^4y [\square_y \Pi(x, y)] h(y) \\ &= \int d^4y \Pi(x, y) \square_y h(y) \quad \text{integrated by parts} \end{aligned} \quad (3.66)$$

therefore,

$$h_0(x) = \int d^4y \Pi(x, y) J(y) \quad (3.67)$$

The equation $\square h_1 = \lambda h_0$ can be rewritten as a new form easy to compute:

$$\begin{aligned} \square_\omega h_1(\omega) &= \lambda h_0^2(\omega) \\ &= \lambda h_0(\omega) h_0(\omega) \\ &= \lambda \int d^4y \Pi(\omega, y) J(y) \int d^4z \Pi(\omega, z) J(z) \end{aligned} \quad (3.68)$$

So the final solution in order of λ is :

$$h(x) = \int d^4y \Pi(x, y) J(y) + \lambda \int d^4y \Pi(\omega, y) J(y) \int d^4z \Pi(\omega, z) J(z) + \mathcal{O}(\lambda^2) \quad (3.69)$$

There is a Feynman diagram to show this process. But we keep this till the later chapter about formal Feynman diagrams and rules. Here, the Feynman rules for this classical field theory examphy are : 1. Draw a point x and a line from x to a new point x_i . 2. Either truncate a line at a source J or let the line branch into two lines adding a new point and a factor of λ . 3. Repeat previous step. 4. The final value for $h(x)$ is given by graphs up to some order in λ with the ends capped by currents $J(x_i)$, the lines replaced by propagators $\Pi(x_i, x_j)$, and all internal points integrated over.

All the above is the content of this chapter, in this chapter, we re-interpret some important contents of the classical field theory which will be helpful to study the quantum field theory later.

Chapter 4

Old-Fashioned Perturbation Theory (OFPT)

In this chapter, we briefly introduce the old-fashioned perturbation theory (OFPT) and then do a little calculation using the theory. At the same time, we recall a bit of confusion that Oppenheimer obtained using OFPT, namely the Lamb shift, which will show us that it is a difficult and important thing to properly handle and clarify all infinite divergences.

One important thing is that there may have some infinities and off-shell cases. And we gonna to try to explain where off-shellness comes from, why you do not need it, but why you want it anyway.

Let's start with the Schwinger's work in OFPT. Just as in quantum mechanics we have leant before, perturbation theory in quantum field theory works by splitting the Hamiltonian up into two parts:

$$H = H_0 + V \quad (4.1)$$

where the H_0 is what we have completely understand and the V can be regarded as small parts¹. We know the state of a system at early times and would like to know the state at late times. For system described by H_0 , we note the complete set as $\{|\phi\rangle\}$, then the eigenequations are:

$$H_0 |\phi\rangle = E |\phi\rangle \quad (4.2)$$

If the energies E are continuous, we shall be able to find an eigenstate $|\psi\rangle$ of the full Hamiltonian with the same energy eigenvalue:

$$H_0 |\psi\rangle = E |\psi\rangle \quad (4.3)$$

Both with the two eigenequations above, we can formally write the linear combination with two states²:

$$|\psi\rangle = |\phi\rangle + \frac{1}{E - H_0} V |\psi\rangle \quad (4.4)$$

¹The difference from quantum mechanics is that in quantum field theory the states often have a continuous range of energies.

²We pick out the state $|\phi\rangle$ that is orthogonal to all other vectors and here we omit the degenerate perturbation problems. Just use the linear algebra.

which is called the *Lippmann-Schwinger equation*³.

Notice that, the algebra part of $\frac{1}{E - H_0}$ is an inverse operator called *Lippmann-Schwinger kernel* which is a kind of *Green's function*:

$$\Pi_{LS} = \frac{1}{E - H_0} \quad (4.5)$$

We define the *transfer matrix* T , satisfying the linear transformation:

$$V |\psi\rangle = T |\phi\rangle \quad (4.6)$$

thus product operator V :

$$T |\phi\rangle = V |\phi\rangle + V \frac{1}{E - H_0} T |\phi\rangle$$

therefore

$$T = V + V \frac{1}{E - H_0} T \quad (4.7)$$

namely

$$T = V + V \Pi_{LS} T \quad (4.8)$$

We expand the equation in order by V :

$$T = V + V \Pi_{LS} V + V \Pi_{LS} V \Pi_{LS} V + \dots \quad (4.9)$$

In Heisenberg's picture, we write the matrix elements of transfer matrix T_{fi} of the initial state mapping to final state⁴:

$$\begin{aligned} \langle \phi_f | T | \phi_i \rangle &= \langle \phi_f | V | \phi_i \rangle + \langle \phi_f | V \Pi_{LS} V | \phi_i \rangle + \dots \\ &= \langle \phi_f | V | \phi_i \rangle + \langle \phi_f | V \frac{1}{E - H_0} | \phi_j \rangle \langle \phi_j | V | \phi_i \rangle + \dots \end{aligned} \quad (4.10)$$

or the matrix product form:

$$T_{fi} = V_{fi} + V_{fi} \Pi_{LS}(j) V_{ji} + V_{fi} \Pi_{LS}(j) V_{jk} \Pi_{LS}(k) V_{ki} + \dots \quad (4.11)$$

where $\Pi_{LS}(k) = \frac{1}{E - E_k}$ and $E = E_i = E_f$ is the energy of the initial and final state we are interested in.

The expansion is the so-called *old-fashioned perturbation theory (OFPT)*.

So how we understand the perturbation's physical picture? Actually, in each term appeared in the equation, creates an intermediate state $|\phi_j\rangle$ which propagates with the propagator $\Pi_{LS}(j)$ until it hits another potential, where it creates a new field $|\phi_k\rangle$ which then propagates and so on, until they hit the final potential factor, which transitions it to the final state. The nice diagrammatic way of drawing this series is the Feynman graphs.

³It is useful in scattering theory. In scattering calculations the V acts at intermediate times to induce transitions among states $|\phi\rangle$ that are assumed to be free (non-interacting) at early and late times. It says the full wavefunction $|\psi\rangle$ is given by the free wavefunction $|\phi\rangle$ plus a scattering term.

⁴insert the complete set and use the Einstein's summation convention

Schwartz introduced the Feynman rule for OFPT and developed it by an important example, an electron scattering off another electron and the Coulomb's law revisited, in his textbook. We don't note this parts rather recalling the Oppenheimer's confusion. We will recall the scattering example till we discuss the Feynman rules in later chapter.

So now, we revisit the second-quantized Hamiltonian from the photon field:

$$H = \int \frac{d^3k}{(2\pi)^3} \omega_k (a_k^\dagger a_k + \frac{1}{2}) \quad (4.12)$$

with $\omega_k = |\vec{k}|$.

The problem is that, if we compute the energy of vacuum⁵, the Hamiltonian gives:

$$E_0 = \langle 0 | H | 0 \rangle = \frac{1}{2} \int \frac{d^3k}{(2\pi)^3} |\vec{k}| = \infty \quad (4.13)$$

which we have used the commutation of the creation and annihilation operator.

That's a bad answer? Fortunately, there is an easy way out of this paradoxical infinity: How do you measure the energy of vacuum? Yes, we cannot. Only energy differences are measurable, and in these differences the *zero-point energy*, the energy of the ground state, drops out. This is the basic idea behind renormalization—infinites can appear in intermediate calculations but they must drop out of physical observables⁶. Because this kind of infinity, the Oppenheimer used to conclude that QED was wrong⁷.

Schwartz used the example of *linear Stark effect* and *Lamb shift* to illustrate the confusions. But we also recall this part in the later chapter which we will develop the formal Feynman diagrams and rules. Let's cut to the chase. What went wrong? In the Stark effect calculation we only had to sum over excited electron states, through $\sum_{m>0} |\psi_m\rangle \langle \psi_m|$, which was finite. For the Lamb shift calculation, the sum was also over photon states, which was divergent. It diverged because the phase space for photons, d^3k , is larger than the suppression, $1/|\vec{k}|$, due to the energies of the intermediate excited states. In terms of Feynman diagrams, the difference is that in the latter case we do not consider interactions with a fixed external field, but integrate over dynamical fields, corresponding to intermediate state photons. Since the photons relevant to the $\langle \psi_0 | H_{int} | \psi_m; 1_k \rangle$ matrix element are the same as the photons relevant to the second, $\langle \psi_m; 1_k | H_{int} | \psi_0 \rangle$ matrix element, the photon lines represent the same state and should be represented by a single line. Thus the diagram contracts, and the Stark effect diagram becomes a loop diagram for the Lamb shift. These pictures are just shorthand for the perturbation expansion. The loop means that there is an unknown momentum, \vec{k} , over which we have to integrate. Momentum must be conserved, but it can split between the atom and the photon in an infinite number of ways.

⁵Even though the energy of vacuum is ∞ , the energy density is limited in local, and the question is that how it affect the curve.

⁶The zero-point energy we have met in cosmology, the Casimir effect.

⁷In fact, the result is not infinite but a finite calculable quantity known as the Lamb shift, which agrees perfectly with data. However, it is instructive to understand Oppenheimer's argument.

There was actually nothing wrong with Oppenheimer's calculation. He did get the answer that OFPT predicts. What he missed was that there are other infinities that eventually cancel this infinity (for example, the electron mass is infinite too, so in fact his conclusion was on the right track). This discussion was really just meant as a preview to demonstrate the complexities we will be up against. To sort out all these infinities, it will be really helpful, but not strictly necessary, to have a formalism that keeps the symmetries, in particular Lorentz invariance, manifest along the way. Although Schwinger was able to tame the infinities using OFPT, his techniques were not for everyone. In his own words, "Like the silicon chips of more recent years, the Feynman diagram was bringing computation to the masses" .

Chapter 5

Cross Sections and Decay Rates

In this chapter, the author provided a useful tool for us to understand the high-energy physics experiments which said that:

The twentieth century witnessed the invention and development of collider physics as an efficient way to determine which particles exist in nature, their properties, and how they interact. In early experiments, such as Rutherford's discovery of the nucleus in 1911 using α -particles or Anderson's discovery of the positron in 1932 from cosmic rays, the colliding particles came from nature. *Colliders provide a great way to study fundamental interactions because they begin with initial states of essentially fixed momenta, i.e. plane waves, and end up with final states, which also have fixed momenta.*

The experimentally measurable quantities that are based on QM are differential probabilities. These probabilities are given by the modulus squared on inner products of states. We can write such inner products as $\langle f; t_f | i; t_i \rangle$, where $|i; t_i\rangle$ is the initial state we start with at time t_i and $\langle f; t_f|$ is the final state we are interested in at some later time t_f . Since quantum field theory is just quantum mechanics with lots of fields, the experimental quantities we will be able to predict are also of the form $|\langle f; t_f | i; t_i \rangle|^2$. One thing should be cleared that the notation of $\langle f; t_f | i; t_i \rangle$ refers to the Schrodinger picture representation, where the states evolve in time. In Heisenberg picture, which will be the default picture for quantum field theory, we leave the states alone and put all the time evolution into an operator. In the special case where we evolve momentum eigenstates from $t = -\infty$ to $t = +\infty$, relevant for collider physics applications, we give the time-evolution operator a special name: **the scattering or S-matrix**. The **S-matrix** is defined as

$$\langle f | S | i \rangle_{\text{Heisenberg}} = \langle f; \infty | i; -\infty \rangle_{\text{Schrodinger}} \quad (5.1)$$

The S-matrix is defined assuming that all of the things that change the state (the interactions) happen in a finite time interval, so that at asymptotic times, $t = \pm\infty$, the states are free of interactions. Free states at $t = \pm\infty$ are known as *asymptotic states*.

But in this chapter, we relate S-matrix elements to scattering cross sections, which are directly measured in collider experiments.

The cross section is a natural quantity to measure experimentally. The cross-sectional area, $\sigma = \pi r^2$, of the nucleus is given by

$$\sigma = \frac{\text{number of particles scattered}}{\text{time} \times \text{number density in beam} \times \text{velocity of beam}} = \frac{1}{T} \frac{1}{\Phi} N \quad (5.2)$$

where the T is the time of experiment and Φ is the incoming flux (exactly the number density \times velocity of beam) In additionally, the number of scattering, N , is determined completely by the *short-distance* interactions among the particles.

The *differential cross section*, $\frac{d\sigma}{d\Omega}$, which gives the number of scattered particles in a certain solid angle $d\Omega$. Classically, this gives us information about the shape of the object or form of the potential off of which the α -particles are scattered¹.

In QM, taking the derivative of both, that says $d\sigma = \frac{1}{T} \frac{1}{\Phi} dN$, while the derivative has the natural meaning in QM, probabilities $P = N/N_{inc}$, where the N is the number of particles scattering into a given area and N_{inc} is the number of incident particles. So the quantum mechanical cross section is then naturally:

$$d\sigma = \frac{1}{T} \frac{1}{\Phi} dP \quad (5.3)$$

where Φ is the flux. But we should notice that the equation is normalized as if the beam has just one particle, and P is now the quantum mechanical probability of scattering. The differential quantities $d\sigma$ and dP are differential in kinematical variables, such as the angles and energies of the final state particles. The differential number of scattering events measures in a collider experiment is;

$$dN = L \times d\sigma \quad (5.4)$$

where L is the *luminosity*, which is defined by this equation.

One more thing, we should to prove the item

$$\int \frac{d^3p}{2\omega_p} \quad (5.5)$$

is Lorentz invariant which is the measure in \mathbb{M}^4 and we have to prove some characters in δ function. At first, there appears a strange term of $1/2\omega_p$. How it appears? That should be a physical problem not a mathematical structure problem. So, the simplest Lorentz invariant in \mathbb{M}^4 (momentum phase space) is

$$\int d^4p \quad (5.6)$$

but we can restrict this from physics so we can change the form.

In consideration of limit of **on-shell condition**, we can add the δ function to restrict it. Notice that

$$\delta^4(p^\mu p_\mu - m^2) = \delta^4(p^0 p_0 - |\vec{p}|^2 - m^2) = \delta^4(p^0 p_0 - \omega_p^2) \quad (5.7)$$

¹This supports a kind of the "microscope" of the shape or potential

where $\omega_p = \sqrt{|\vec{p}|^2 + m^2}$ or $E_p = \sqrt{|\vec{p}|^2 + m^2}$.

Also we have to restrict the energy. Because the $\omega_p = E_p = \sqrt{|\vec{p}|^2 + m^2}$, so we have to add the condition by $\theta(p^0)$ function to ensure the **positive solution**. Thus, the Lorentz invariant form is

$$\int d^4p \delta^4(p^0 p_0 - \omega_p^2) \theta(p^0) \quad (5.8)$$

Now, we introduce the character of δ function,

$$\delta(g(x)) = \sum_i \frac{\delta(x - x_i)}{g'(x_i)} \quad \text{where zero point } g(x_i) = 0 \quad (5.9)$$

Therefore, the zero points of $p^0 p_0 - \omega_p^2$ are $p_1^0 = \omega_p$ and $p_2^0 = -\omega_p$, so the $\delta(p^0 p_0 - \omega_p^2)$ can be written as

$$\delta^4(p^0 p_0 - \omega_p^2) = \left[\frac{\delta^4(p_0 - \omega_p)}{2\omega_p} - \frac{\delta^4(p_0 + \omega_p)}{2\omega_p} \right] \quad (5.10)$$

where the second item implies the negative energy solution.

So, the computation :

$$\begin{aligned} \int d^4p \delta^4(p^0 p_0 - \omega_p^2) \theta(p_0) &= \int d^4p \left[\frac{\delta^4(p_0 - \omega_p)}{2\omega_p} - \frac{\delta^4(p_0 + \omega_p)}{2\omega_p} \right] \theta(p_0) \\ &= \int dp_0 \int d^3p \left[\frac{\delta^4(p_0 - \omega_p)}{2\omega_p} - \frac{\delta^4(p_0 + \omega_p)}{2\omega_p} \right] \theta(p_0) \\ &= \int dp_0 \int d^3p \frac{\delta(p_0 - \omega_p)}{2\omega_p} \theta(p_0) \\ &= \int \frac{d^3p}{2\omega_p} \end{aligned} \quad (5.11)$$

Therefore,

$$\int \frac{d^3p}{2\omega_p} \quad (5.12)$$

is Lorentz invariant.

Now, let's turn back to our main topic.

to be supplemented

Chapter 6

The S-Matrix and Time-Ordered Products

In this chapter here, our main topic is the LSZ (Lehmann-Symanzik-Zimmermann) reduction formula. As discussed in Chapter 5, scattering experiments have been an efficient way to determine the particles that exist and how they interact. As we know, all of the interesting interacting physics is encoded in how often given initial states produce given final states, that is, in the *S-matrix*. We assume that the process happened in finite time which fitting to the collider scattering experiments. This also means that if there were always interactions, it would not be possible to set up our initial states at $t = -\infty$ or find the desired final states at $t = +\infty$. Without interactions at asymptotic times, the states we scatter can be defined as on-shell one-particle states of given momenta, known as **asymptotic states**. We gonna to derive an expression for the S-matrix using only that the system is free at asymptotic times. And after that, we will derive and use the Feynman rules which is a powerful weapon to the interacting theory.

Before the driving, we give the LSZ reduction formula at first. The LSZ relates the S-matrix elements $\langle f|S|i\rangle$ for n asymptotic momentum eigenstates to an expression involving the quantum fields $\phi(x)$:

$$\langle f|S|i\rangle = \left[i \int d^4x_1 e^{-ip_1x_1} (\square + m^2) \right] \cdots \left[i \int d^4x_n e^{ip_nx_n} \right] \times \langle \Omega | T \{ \phi(x_1) \phi(x_2) \phi(x_3) \cdots \phi(x_n) \} | \Omega \rangle$$

(6.1)

to be supplemented

Chapter 7

Feynman Rules

We have already seen that the scattering cross sections are naturally expressed in terms of time-ordered products of fields. The S-matrix has the form

$$\langle f | S | i \rangle \sim \langle \Omega | T \{ \phi(x_1) \cdots \phi(x_n) \} | \Omega \rangle \quad (7.1)$$

where $|\Omega\rangle$ is the ground state or vacuum in the interacting theory. In this expression the fields $\phi(x)$ are not free but are the full interacting quantum fields. We also know that the time-ordered product of two fields is given by the Feynman propagator in the free theory:

$$D_F(x, y) \equiv \langle 0 | T \{ \phi_0(x) \phi_0(y) \} | 0 \rangle = \lim_{\epsilon \rightarrow 0} \int \frac{d^4 k}{(2\pi)^4} \frac{i}{k^2 - m^2 + i\epsilon} e^{ik(x-y)} \quad (7.2)$$

where $|0\rangle$ is the ground state in the free theory.

There are position-space Feynman rules, for calculating time-ordered products, and also momentum-space Feynman rules, for calculating S-matrix elements. The momentum-space Feynman rules are by far the more important- they provide an extremely efficient way to set up calculations of physical results in quantum field theory. The momentum-space Feynman rules are the main result of this Part I.

At first, we will derive the Feynman rules from the Lagrangian formulation of time evolution and quantization which is the quickest way to connect Feynman diagrams to classical field theory. Then, we will derive it based on expansion of the full interacting Hamiltonian around the free Hamiltonian.

In chapter 2, we got the free field theory which implies two top priority formula:

$$[\phi(\vec{x}, t), \phi(\vec{x}', t)] = 0 \quad (7.3)$$

$$[\phi(\vec{x}, t), \partial_t \phi(\vec{x}', t)] = i\hbar \delta^3(\vec{x} - \vec{x}') \quad (7.4)$$

which called the equal-time commutation relation.

However, we used to set that $\hbar = c = 1$ in previous chapters. The \hbar here in (7.4) just imply the different from classical limit. And we used to derive the free field equation $(\square + m^2)\phi = 0$ by the Euler-Lagrange equation from field theory. While in an interacting theory, we must

generalize these equations to specify how the dynamics is determined. Then we quantize the Heisenberg motion equation from quantum mechanics

$$i\frac{\partial}{\partial t}\phi(x) = [\phi, H] \quad (7.5)$$

which we regard it as one natural approach for an interacting quantum field theory lead to the Hamiltonian derivation of the Feynman rules in the later content.

Our first topic will be developed by the Lagrange approach where the Hamilton's equations are replaced by the Euler-Lagrange equations (we will simply the note by "E-L equations") derived from a Lagrangian \mathcal{L} and the generalization of $(\square + m^2)\phi = 0$. We also assume (7.3) and (7.4) are still satisfied¹. This is a natural assumption, since the Hilbert space for the interacting theory is the same as the free theory. I think this is very physical which we can understand from physics. The equal-time commutation relation of (7.3) implies that the **causality** is satisfied in quantum field theory. Because at the same time but at different points in space, all operators (or say physical quantity), in particular fields, should be simultaneously observable and commute otherwise our causality will broke up.

Now, let us talk about the (7.4). This formula can be seen as QFT version of the basic commutation relation in quantum mechanics of $[\hat{x}, \hat{p}] = i\hbar$. We can understand the $\phi(x)$ is the QFT version of x since we can know that $\phi_0(x)|0\rangle \sim |x\rangle$ and the \hat{p} is also. In another word, the (7.4) is the equivalent of the canonical commutation relation from the quantum mechanics $[\hat{x}, \hat{p}] = i\hbar$. It indicates that a quantity and its time derivative are not simultaneously observable, the hallmark of the **uncertainty principle**.

We have known how to calculate $\langle 0|T\{\phi(x)\phi(x')\}|0\rangle$ in the free theory. Now we calculate this commutator in an interacting theory. Firstly, we prove the formula below:

$$(\square + m^2) \langle \Omega|T\{\phi(x)\phi(x')\}|\Omega\rangle = \langle \Omega|T\{(\square + m^2)\phi(x)\phi(x')\}|\Omega\rangle - i\hbar\delta^4(x - x')$$

where the $|\Omega\rangle$ is the vacuum in the interacting theory. The $\delta^4(x - x')$ on right side is critically important. It signifies the difference between the classical and quantum theories in a way that will be clear shortly. Notice that:

$$\begin{aligned} \partial_t \langle \Omega|T\{\phi(x)\phi(x')\}|\Omega\rangle &= \partial_t [\langle \Omega|\phi(x)\phi(x')|\Omega\rangle \theta(t - t') + \langle \Omega|\phi(x')\phi(x)|\Omega\rangle \theta(t' - t)] \\ &= \langle \Omega|\partial_t\phi(x)\phi(x')|\Omega\rangle \theta(t - t') + \langle \Omega|\phi(x')\partial_t\phi(x)|\Omega\rangle \theta(t' - t) \\ &\quad + \langle \Omega|\phi(x)\phi(x')|\Omega\rangle \partial_t\theta(t - t') + \langle \Omega|\phi(x')\phi(x)|\Omega\rangle \partial_t\theta(t' - t) \\ &= \langle \Omega|T\{\partial_t\phi(x)\phi(x')\}|\Omega\rangle + \langle \Omega|\phi(x)\phi(x')|\Omega\rangle \delta(t - t') - \langle \Omega|\phi(x')\phi(x)|\Omega\rangle \delta(t - t') \\ &= \langle \Omega|T\{\partial_t\phi(x)\phi(x')\}|\Omega\rangle + \langle \Omega|[\phi(x)\phi(x') - \phi(x')\phi(x)]|\Omega\rangle \delta(t - t') \\ &= \langle \Omega|T\{\partial_t\phi(x)\phi(x')\}|\Omega\rangle + \langle \Omega|[\phi(x), \phi(x')]\delta(t - t') \\ &= \langle \Omega|T\{\partial_t\phi(x)\phi(x')\}|\Omega\rangle + 0 \\ &= \langle \Omega|T\{\partial_t\phi(x)\phi(x')\}|\Omega\rangle \end{aligned} \quad (7.6)$$

¹We have do proof of this through GP. Zhang's GFT course

and we used the $\partial_x \theta(x) = \delta(x)$ and the $\delta(t - t')$ forces the $t = t'$ therefore there exists the equal time and we use the equal time commutator $[\phi(\vec{x}, t), \phi(\vec{x}', t)] = 0$ leading the vanish of the second term from above. Simply,

$$\partial_t \langle \Omega | T \{ \phi(x) \phi(x') \} | \Omega \rangle = \langle \Omega | T \{ \partial_t \phi(x) \phi(x') \} | \Omega \rangle \quad (7.7)$$

Take the second derivative:

$$\begin{aligned} \partial_t^2 \langle \Omega | T \{ \phi(x) \phi(x') \} | \Omega \rangle &= \partial_t \langle \Omega | T \{ \partial_t \phi(x) \phi(x') \} | \Omega \rangle \\ &= \partial_t [\langle \Omega | \partial_t \phi(x) \phi(x') | \Omega \rangle \theta(t - t') + \langle \Omega | \phi(x') \partial_t \phi(x) | \Omega \rangle \theta(t' - t)] \\ &= \langle \Omega | T \{ \partial_t^2 \phi(x) \phi(x') \} | \Omega \rangle + \langle \Omega | \partial_t \phi(x) \phi(x') - \phi(x') \partial_t \phi(x) | \Omega \rangle \delta(t - t') \\ &= \langle \Omega | T \{ \partial_t^2 \phi(x) \phi(x') \} | \Omega \rangle + \langle \Omega | [\partial_t \phi(x), \phi(x')] | \Omega \rangle \delta(t - t') \\ &= \langle \Omega | T \{ \partial_t^2 \phi(x) \phi(x') \} | \Omega \rangle - i\hbar \delta^3(\vec{x} - \vec{x}') \delta(t - t') \\ &= \langle \Omega | T \{ \partial_t^2 \phi(x) \phi(x') \} | \Omega \rangle - i\hbar \delta^4(x - x') \end{aligned} \quad (7.8)$$

where we have used $\partial_t \phi(x') = 0$ and the $\delta(t - t')$ forces the equal time again and we can use the equal-time commutator $[\phi(x), \partial_t \phi(x')] = i\hbar \delta^3(\vec{x} - \vec{x}')$ and the normalization of $\langle \Omega | \Omega \rangle = 1$.

Here, we just get the relation of $\partial^2 / \partial t^2$ but our goal is $\square + m^2 = \partial_t^2 - \partial_x^2 + m^2$. The m^2 can be seen as a constant, so we can insert the m^2 in anywhere. Then we have to deal with the second derivative of space ∂_x^2 or ∇^2 . Same as the above, we have

$$\begin{aligned} \partial_x \langle \Omega | T \{ \phi(x) \phi(x') \} | \Omega \rangle &= \partial_x [\langle \Omega | \phi(x) \phi(x') | \Omega \rangle \theta(t - t') + \langle \Omega | \phi(x') \phi(x) | \Omega \rangle \theta(t' - t)] \\ &= \langle \Omega | \partial_x^2 \phi(x) \phi(x') | \Omega \rangle \theta(t - t') + \langle \Omega | \phi(x') \partial_x^2 \phi(x) | \Omega \rangle \delta(t' - t) \\ &= \langle \Omega | T \{ \partial_x^2 \phi(x) \phi(x') \} | \Omega \rangle \end{aligned} \quad (7.9)$$

where we have used the formula $\partial_x \phi(x') = 0$ and $\partial_x \theta(t - t') = 0$.

Above all, we get that

$$(\square + m^2) \langle \Omega | T \{ \phi(x) \phi(x') \} | \Omega \rangle = \langle \Omega | T \{ (\square + m^2) \phi(x) \phi(x') \} | \Omega \rangle - i\hbar \delta^4(x - x') \quad (7.10)$$

For example, in the free theory, $\langle 0 | T \{ \phi_0(x) \phi_0(y) \} | 0 \rangle = \hbar D_F(x, y)$ and $(\square + m^2) \phi_0(X) = 0$, then (7.10) implies

$$(\square_x + m^2) D_F(x, y) = -\delta^4(x - y) \quad (7.11)$$

To simplify the equation expression, we note $\langle \dots \rangle = \langle \Omega | T \{ \dots \} | \Omega \rangle$ for time-ordered correlation functions in the interacting theory, (7.10) can be written as

$$(\square + m^2) \langle \phi(x) \phi(x') \rangle = \langle (\square + m^2) \phi(x) \phi(x') \rangle - i\hbar \delta^4(x - x') \quad (7.12)$$

Let's insert more fields which can be expressed as a normal form

$$\square_x \langle \phi(x) \phi(x_1) \cdots \phi(x_n) \rangle = \langle \square_x \phi(x) \phi(x_1) \cdots \phi(x_n) \rangle - i\hbar \sum_j \delta^4(x - x_j) \langle \phi(x_1) \cdots \phi(x_{j-1}) \phi(x_{j+1}) \cdots \phi(x_n) \rangle$$

It is not difficult to imagine that the above formula should be summarized from some law. So we calculate the 3 fields $\square_x \langle \phi(x) \phi(x_1) \phi(x_2) \rangle$ which the M.D. Schwartz recommend we

calculating it. Then, we can summarize the law from the 3 fields and our original 2 field formula. Similar to the 2 field formula, the $\square_x = \partial_t^2 - \partial_x^2$, is essential to our calculating especially the time second derivative which will calculate on the $\theta(t)$ part while the ∂_x does not.

For the time-order product of 3 fields, they will appear the permutation problem. So we simply note the time first derivative

$$\begin{aligned} \partial_t \langle \phi(x)\phi(x_1)\phi(x_2) \rangle &= \partial_t \langle \Omega | T\{\phi(x)\phi(x_1)\phi(x_2)\} | \Omega \rangle \\ &= \sum_{\text{Permutation}}^{A_3^3=6} \partial_t [\langle \Omega | \phi(x)\phi(x_1)\phi(x_2) | \Omega \rangle \theta(t-t_1)\theta(t_1-t_2)] \end{aligned} \quad (7.13)$$

Now, focus on the single kind " $\phi(x)\phi(x_1)\phi(x_2)$ ":

$$\begin{aligned} &\partial_t [\langle \Omega | \phi(x)\phi(x_1)\phi(x_2) | \Omega \rangle \theta(t-t_1)\theta(t_1-t_2)] \\ &= \langle \Omega | \partial_t \phi(x)\phi(x_1)\phi(x_2) | \Omega \rangle \theta(t-t_1)\theta(t_1-t_2) + \langle \Omega | \phi(x)\phi(x_1)\phi(x_2) | \Omega \rangle \delta(t-t_1)\theta(t_1-t_2) \end{aligned} \quad (7.14)$$

which we need to emphasize that the derivative of t or x only take effect on the quantity involve the t and x . Any permutations and combinations both have the similar form as above:

$$\langle \Omega | \partial_t \phi(x)\phi(x_1)\phi(x_2) | \Omega \rangle \theta(t-t_1)\theta(t_1-t_2) \quad (P.1)$$

$$\langle \Omega | \partial_t \phi(x)\phi(x_2)\phi(x_1) | \Omega \rangle \theta(t-t_2)\theta(t_2-t_1) \quad (P.2)$$

$$\langle \Omega | \phi(x_1)\partial_t \phi(x)\phi(x_2) | \Omega \rangle \theta(t_1-t)\theta(t-t_2) \quad (P.3)$$

$$\langle \Omega | \phi(x_1)\phi(x_2)\partial_t \phi(x) | \Omega \rangle \theta(t_1-t_2)\theta(t_2-t) \quad (P.4)$$

$$\langle \Omega | \phi(x_2)\partial_t \phi(x)\phi(x_1) | \Omega \rangle \theta(t_2-t)\theta(t-t_1) \quad (P.5)$$

$$\langle \Omega | \phi(x_2)\phi(x_1)\partial_t \phi(x) | \Omega \rangle \theta(t_2-t_1)\theta(t_1-t) \quad (P.6)$$

Notice that, the derivative of x in this formula gives $\partial_x \phi(x_j) = 0$ and $\partial_x \theta(t) = 0$. This means that the derivative of x can be inserted in anywhere in the formula. And the ∂_x^2 is so. We can directly insert the ∂_x^2 in it while ∂_t^2 does not. Then we can say: the d'Alembert operator can be written as the follow:

$$\square_x \langle \phi(x)\phi(x_1)\phi(x_2) \rangle = \langle \square_x \phi(x)\phi(x_1)\phi(x_2) \rangle + \text{Remainder}\{t\} \quad (7.15)$$

which the $\text{Remainder}\{t\}$ implies that the remainder is caused by time derivative.

Since the \square_x in the second derivative especially the ∂_t^2 , we can decompose the remainder as

$$\text{Remainder}\{t\} = \partial_t (\text{first remainder}) + (\text{second remainder}) \quad (7.16)$$

from the time first derivative and the second derivative.

Let's check the time first derivative. We can easily write the first remainder by taking the time derivative of (P.1)~(P.6). Then, as for the ∂_t , we have

$$\langle \Omega | \phi(x)\phi(x_1)\phi(x_2) | \Omega \rangle \delta(t-t_1)\theta(t_1-t_2) \quad (F.1)$$

$$\langle \Omega | \phi(x) \phi(x_2) \phi(x_1) | \Omega \rangle \delta(t - t_2) \theta(t_2 - t_1) \quad (\text{F.2})$$

$$\langle \Omega | \phi(x_1) \phi(x) \phi(x_2) | \Omega \rangle \delta(t - t_2) \theta(t_1 - t) \quad (\text{F.3})$$

$$\langle \Omega | \phi(x_2) \phi(x) \phi(x_1) | \Omega \rangle \delta(t - t_1) \theta(t_2 - t) \quad (\text{F.4})$$

$$- \langle \Omega | \phi(x_1) \phi(x_2) \phi(x) | \Omega \rangle \delta(t - t_2) \theta(t_1 - t_2) \quad (\text{F.5})$$

$$- \langle \Omega | \phi(x_2) \phi(x_1) \phi(x) | \Omega \rangle \delta(t - t_1) \theta(t_2 - t_1) \quad (\text{F.6})$$

$$- \langle \Omega | \phi(x_1) \phi(x) \phi(x_2) | \Omega \rangle \delta(t - t_1) \theta(t - t_2) \quad (\text{F.7})$$

$$- \langle \Omega | \phi(x_2) \phi(x) \phi(x_1) | \Omega \rangle \delta(t - t_2) \theta(t - t_1) \quad (\text{F.8})$$

which we have used the $\partial_x \theta(x) = \delta(x)$ and we order all t be front in the δ function legally.

Samely, the remainder caused by ∂_t^2 can be listed as followings (just replace the $\phi(x)$ to $\partial_t \phi(x)$):

$$\langle \Omega | \partial_t \phi(x) \phi(x_1) \phi(x_2) | \Omega \rangle \delta(t - t_1) \theta(t_1 - t_2) \quad (\text{S.1})$$

$$\langle \Omega | \partial_t \phi(x) \phi(x_2) \phi(x_1) | \Omega \rangle \delta(t - t_2) \theta(t_2 - t_1) \quad (\text{S.2})$$

$$\langle \Omega | \phi(x_1) \partial_t \phi(x) \phi(x_2) | \Omega \rangle \delta(t - t_2) \theta(t_1 - t) \quad (\text{S.3})$$

$$\langle \Omega | \phi(x_2) \partial_t \phi(x) \phi(x_1) | \Omega \rangle \delta(t - t_1) \theta(t_2 - t) \quad (\text{S.4})$$

$$- \langle \Omega | \phi(x_1) \phi(x_2) \partial_t \phi(x) | \Omega \rangle \delta(t - t_2) \theta(t_1 - t_2) \quad (\text{S.5})$$

$$- \langle \Omega | \phi(x_2) \phi(x_1) \partial_t \phi(x) | \Omega \rangle \delta(t - t_1) \theta(t_2 - t_1) \quad (\text{S.6})$$

$$- \langle \Omega | \phi(x_1) \partial_t \phi(x) \phi(x_2) | \Omega \rangle \delta(t - t_1) \theta(t - t_2) \quad (\text{S.7})$$

$$- \langle \Omega | \phi(x_2) \partial_t \phi(x) \phi(x_1) | \Omega \rangle \delta(t - t_2) \theta(t - t_1) \quad (\text{S.8})$$

Now, recall the method we used in the 2 fields problem. We constructed the equal-time commutator to simplify the formula and remove the awful terms. Now, we have eight equations from (F.1) to (F.8) and eight equations from (S.1) to (S.8). We try to constuct the commutator. To make this easier, we omit the $|\Omega\rangle$ and $\langle\Omega|$ symbols and use two three-dimensional vector symbols to represent that they are equal in time like that $\phi(\vec{x})$ and $\phi(\vec{x}_1)$ are equal time.

We observe the (F.1)~(F.8) and we can find that the $\delta(t - t_1)$ enforce the equal time $t = t_1$. Then we can combine the (F.1), (F.4), (F.6) and (F.7) and give

$$\begin{aligned} & \phi(\vec{x}) \phi(\vec{x}_1) \phi(x_2) \theta(t - t_2) - \phi(\vec{x}_1) \phi(\vec{x}) \phi(x_2) \theta(t - t_2) + \phi(x_2) \phi(\vec{x}) \phi(\vec{x}_1) \theta(t_2 - t) - \phi(x_2) \phi(\vec{x}_1) \phi(\vec{x}) \theta(t_2 - t) \\ &= [\phi(\vec{x}), \phi(\vec{x}_1)] \phi(x_2) \theta(t - t_2) + \phi(x_2) [\phi(\vec{x}), \phi(\vec{x}_1)] \theta(t_2 - t) \end{aligned} \quad (\text{7.17})$$

equally, the $\delta(t - t_2)$ enforces the equal time $t = t_2$, and we have

$$\begin{aligned} & \phi(\vec{x}) \phi(\vec{x}_2) \phi(x_1) \theta(t - t_1) - \phi(\vec{x}_2) \phi(\vec{x}) \phi(x_1) \theta(t - t_1) + \phi(x_1) \phi(\vec{x}) \phi(\vec{x}_2) \theta(t_1 - t) - \phi(x_1) \phi(\vec{x}_2) \phi(\vec{x}) \theta(t_1 - t) \\ &= [\phi(\vec{x}), \phi(\vec{x}_2)] \phi(x_1) \theta(t - t_1) + \phi(x_1) [\phi(\vec{x}), \phi(\vec{x}_2)] \theta(t_1 - t) \end{aligned} \quad (\text{7.18})$$

while the equal-time commutators

$$\begin{aligned} [\phi(\vec{x}), \phi(\vec{x}_1)] &= [\phi(\vec{x}, t), \phi(\vec{x}_1, t)] = 0 \\ [\phi(\vec{x}), \phi(\vec{x}_2)] &= [\phi(\vec{x}, t), \phi(\vec{x}_2, t)] = 0 \end{aligned} \quad (7.19)$$

imply

$$\text{first remainder} = 0 \quad (7.20)$$

namely

$$\partial_t(\text{first remainder}) = 0 \quad (7.21)$$

This proves that our mind of constructing the commutator is success. Then we now deal with the second remainder. Equally, we just replace the $\phi(x)$ to $\partial_t \phi(x)$ as we already seen from (S.1) to (S.8) and replace the commutator relation to $[\phi(\vec{x}, t), \partial_t(\vec{x}', t)] = i\hbar \delta^3(\vec{x} - \vec{x}')$. We give back the $\delta(t - t_1)$ and $\delta(t - t_2)$ and re-write the (7.17) and (7.18). Then, we get

$$\begin{aligned} & [\phi(\vec{x}), \phi(\vec{x}_1)] \phi(x_2) \delta(t - t_1) \theta(t - t_2) + \phi(x_2) [\phi(\vec{x}), \phi(\vec{x}_1)] \delta(t - t_1) \theta(t_2 - t) \\ &= -i\hbar \delta^3(\vec{x} - \vec{x}_1) \delta(t - t_1) \phi(x_2) \theta(t - t_2) - i\hbar \delta^3(\vec{x} - \vec{x}_1) \phi(x_2) \theta(t_2 - t) \\ &= -i\hbar \delta^4(x - x_1) [\phi(x_2) \theta(t - t_2) + \phi(x_2) \theta(t_2 - t)] \end{aligned} \quad (7.22)$$

and

$$\begin{aligned} & [\phi(\vec{x}), \phi(\vec{x}_2)] \phi(x_1) \delta(t - t_2) \theta(t - t_1) + \phi(x_1) [\phi(\vec{x}), \phi(\vec{x}_2)] \delta(t - t_2) \theta(t_1 - t) \\ &= -i\hbar \delta^4(x - x_2) [\phi(x_1) \theta(t - t_1) + \phi(x_1) \theta(t_1 - t)] \end{aligned} \quad (7.23)$$

They show us that the second remainder have such a construction

$$\text{Remainder}\{t\} = \text{second remainder} = -i\hbar \sum_j \delta^4(x - x_j) [\text{time ordered product, no } \phi(x) \text{ and no } \phi(x_j)] \quad (7.24)$$

However, we just solve the 3 fields' $\square_x \langle \phi(x) \phi(x_1) \phi(x_2) \rangle$. The time ordered product just implies one field's time order. But, let's recall the 2 field formula of (7.10) then we can summarize the n-fields' formula from the law of consturction of communtators:

$$\square_x \langle \phi(x) \phi(x_1) \cdots \phi(x_n) \rangle = \langle \square_x \phi(x) \phi(x_1) \cdots \phi(x_n) \rangle - i\hbar \sum_j \delta^4(x - x_j) \langle \phi(x_1) \cdots \phi(x_{j-1}) \phi(x_{j+1}) \cdots \phi(x_n) \rangle \quad (7.25)$$

which is useful!

We assunme that the quantum field satisfies the same equations of motion as the classical field we developed in Chapter 3. In particular, if our Lagrangian has the form

$$\mathcal{L} = -\frac{1}{2} \phi(\square + m^2) \phi + \mathcal{L}_{\text{int}}[\phi] \quad (7.26)$$

then the motion equation developed from E-L equation is

$$(\square + m^2) \phi - \mathcal{L}'_{\text{int}}[\phi] = 0 \quad (7.27)$$

where $\mathcal{L}_{\text{int}}[\phi] = d\mathcal{L}_{\text{int}}[\phi]/d\phi$.

Then, we can get the following with (7.25):

$$\boxed{(\square_x + m^2) \langle \phi_x \phi_1 \cdots \phi_n \rangle = \langle \mathcal{L}'_{\text{int}}[\phi_x] \phi_1 \cdots \phi_n \rangle - i\hbar \sum_j \delta^4(x - x_j) \langle \phi_1 \cdots \phi_{j-1} \phi_{j+1} \cdots \phi_n \rangle} \quad (7.28)$$

where $\phi_x = \phi(x)$ and $\phi_j = \phi(x_j)$.

The (7.28) is called **Schwinger-Dyson equations**.

The Schwinger–Dyson equations encode the difference between the classical and quantum theories. Note that their derivation did not require any specification of the dynamics of the theory, only that the canonical commutation relations in (7.19) are satisfied. In particular, in a classical theory, $[\phi(x', t), \partial_t \phi(x, t)] = 0$ and therefore classical timeordered correlation functions would satisfy a similar equation but without the $\delta^4(x - x_j)$ terms (i.e. $\hbar = 0$). That is, in a classical theory, correlation functions satisfy the same differential equations as the fields within the correlation functions. In a quantum theory, that is true only up to δ -functions, which in this context are also called contact interactions. These contact interactions allow virtual particles to be created and destroyed, which permits closed loops to form in the Feynman diagrammatic expansion, as we will now see.

to be supplemented. It's Position-Space Feynman Rules