

Lectures on Stochastic Quantization of Φ_3^4

Part I: From QM to EQFT

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This is an extended version of the material presented in the first part of the course (10h) where I give an overview of the mathematical structures underlying quantum mechanics and the modifications needed to deal with special relativity, all with an emphasis on the Euclidean approach and probabilistic constructions. The second part will take from here to develop the stochastic quantisation of a particular model, the Φ_3^4 Euclidean quantum field theory.

This script has not been seriously revised and can contain typos and inconsistencies, use at your own risk. Also I didn't try to cite all the literature associated to the presented material. Where useful I cite sources where the reader can find more detailed discussions.

These notes have been written with $\text{T}_{\text{E}}\text{X}_{\text{M}}\text{A}^{\text{C}}\text{S}$.

1 Introduction

The aim of these lectures is to illustrate the idea of *stochastic quantization* using as motivating example the construction of the Φ_3^4 Euclidean quantum field theory (EQFT).

(Bosonic) EQFTs are certain (complicated and not very explicit) probability measures μ on spaces of Schwarz distributions $\mathcal{S}'(\mathbb{R}^d)$ over the Euclidean space \mathbb{R}^d which have applications in mathematical physics, in particular in the construction of models of relativistic quantum mechanics, the so called quantum field theories (QFTs).

The stochastic quantization idea addresses the problem of EQFT via the construction of an auxiliary map F_μ which push forward a Gaussian measure γ (possibly defined on a very different space) into the measure μ describing the EQFT: $\mu = F_\mu \gamma$. The measure γ should be thought as a convenient source of randomness and will be chosen such as to facilitate the analysis: Gaussian measures are very well understood and several tools are available for them. The map F_μ on the other hand is expected to contain all the difficulties. However it turns out that tools and ideas from the analysis of PDE or even new analysis tools (like Hairer's regularity structures) can be put forward to obtain interesting results.

The leitmotif of these lectures will be to explore the possibility of obtaining as much information as possible on μ via the stochastic quantization map F_μ . It is therefore useful to fix our *partito preso* in a very general definition which seems to fit all the cases I care about:

Definition 1. A stochastic quantization of a probability measure μ is a pair (F_μ, W) of a map F_μ and a Gaussian r.v. W (on a given probability space) such that the r.v.

$$\phi = F_\mu(W)$$

has law μ .

There are various possibilities for the choice of F_μ , some of them quite equivalent in their usefulness. During these lectures we will concentrate on a specific measure μ , the Φ_3^4 Euclidean quantum field theory and on a specific way to construct F_μ first suggested in this context by Parisi and Wu: introducing a fictitious time and considering a (over-damped) Langevin dynamics driven by a space–time white noise and with invariant measure μ .

Part of my goal is to introduce all these concepts.

For the moment I would like to concentrate on *motivating* why we are interested in such measures μ in the first place, i.e. what an EQFT has to do with quantum mechanics. The link between quantum mechanics and probability theory (and statistical mechanics) has been discovered and exploited in the '60-'70 by people like Segal, Nelson, Symanzik, Simon, Jaffe, Glimm, Osterwalder, Schrader, Rosen, Guerra, Gallavotti, Jona-Lasinio, Fröhlich, ... and many others. It is not my aim here to provide a full account of the development of these ideas, the interested reader can refer e.g. to the book of Glimm and Jaffe (James Glimm and Arthur Jaffe, *Quantum Physics: A Functional Integral Point of View*, 2nd ed. (New York: Springer-Verlag, 1987), [//www.springer.com/gb/book/9780387964775](http://www.springer.com/gb/book/9780387964775)).

I want to introduce the basic mathematical structure of quantum mechanics and state the basic problems related to it: the construction of suitable representations of the algebra of observables with interesting dynamics. Very intuitively this comes down to solving certain ODE in non-commutative variables, the non-commutativity is imposed by the indetermination principle of Heisenberg (in one form or another). When trying to construct a quantum theory which satisfies also the requirements of special relativity (i.e. no physical influence can travel faster than light, covariance wrt. the Poincaré group of Minkowski space) one need to introduce the concepts of observables localized in certain space-time regions and impose constraints on space-like-separated observables. The so called Haag-Kastler axioms of local QFT formalize the basic structures one would like to obtain from these algebras of observables.

Constructing interesting examples of these algebras, especially for physical geometries, has revealed itself a very difficult task and naive attempts (i.e. perturbation theory around known solutions) run into serious problems of divergences: some of the correction to the zero-order behavior are ill-defined and formally infinitely-large. The main problems are related to the local non-commutativity of the algebras which “creates” such singular behavior. In order to cope with these difficulties it has been proposed, first by the physicist, then rigorously by Nelson and Osterwalder-Schräder, to consider the Euclidean “shadow” of the Minkowski theory, by performing an analytic continuation in the time variable $t \in \mathbb{R}$ to the imaginary axis $ix_0 \in i\mathbb{R}$. This transformation sends the Minkowski metric $x_1^2 + \dots + x_n^2 - t^2$ of the n -dimensional Minkowski space \mathbb{M}^n into the Euclidean metric $x_1^2 + \dots + x_n^2 + x_0^2$ of the space \mathbb{R}^{n+1} . Technically this analytic continuation is performed at the level of Wightman functions: a system of correlation functions which allows the reconstruction of an algebra of local observables. The analytically continued Wightman functions become Schwinger functions. It turns out (a bit surprisingly) that Schwinger functions can be constructed via probability theory: they can be taken to be the moments of a measure μ on $\mathcal{S}'(\mathbb{R}^d)$ with $d = n + 1$. In order to allow to go back from Schwinger functions to Wightman functions and to the algebra of local observables one need to verify a series of conditions, i.e. axioms for Schwinger functions. The key conditions are reflection positivity, Euclidean covariance and certain moment estimates. Additionally one would like from the model to satisfy additional conditions but in these lectures we will not study them and concentrate on these three. For the scope of these lecture we can therefore preliminary define

Definition 2. An EQFT μ is a probability measure on $\mathcal{S}'(\mathbb{R}^d)$ satisfying: reflection positivity, Euclidean invariance and certain moment estimates (to be specified later on).

We need to define all the concepts entering this definition: this will be the goal of the first few lectures. The less natural is reflection positivity (RP) which is a however the key property which encodes the quantum theory within a probabilistic object. RP seems innocuous, but in connection with Euclidean invariance give hard constraints for μ to be an EQFT.

The main theorem we would like to prove in these lectures is

Theorem 3. There exists a (two-parameter) family $\{\mu\}$ of EQFT on $\mathcal{S}'(\mathbb{R}^3)$. We call them the Φ_3^4 model.

The name of the model contains two numbers: 3 refers to the dimension of \mathbb{R}^3 while the 4 is a reminder of the form of the approximations we will use to show Theorem 3.

A similar statement (and more general ones) are true in two (Euclidean) dimension, i.e. the Φ_2^4 model. We concentrate in dimension three since the proof is more difficult and showcase certain features of the stochastic quantization approach. Sometimes along the lectures we will refer to the Φ_2^4 for illustrative purpose.

The literature on the Euclidean approach is vast and again I will not make any attempt to present it completely, however it mainly exploits techniques and points of view others than the stochastic quantization approach (be it phase cell expansion, renormalization group, ...). A lot is known of EQFT even beyond the Φ_d^4 models and many research produced profound and technically difficult pieces of work. Later on I might spend some time reviewing some the highlights.

The rigorous construction of Φ_3^4 via stochastic quantization has been developed very recently and I will try to mention all the relevant references so that the interested reader can have an overview of the current state of this line of research.

Let me go back for the moment to the Definition 1 of SQ and to its spirit. As I already mentioned, the idea is to use a convenient “universal” source of randomness to construct a random field with given law. This broad definition fits also to Ito's approach to the construction of diffusion processes with prescribed characteristics. Indeed consider the (strong) solution X of a stochastic differential equation (SDE) of the form

$$dX_t = b(X_t)dt + \sigma(X_t)dW_t, \quad X_0 = x$$

on \mathbb{R}^k where W is a (finite-dimensional) Brownian motion and b, σ are prescribed vector-fields on \mathbb{R}^k . This strong solution is really a map from W to $X = F(W)$ and Ito's motivation to introduce SDEs was mainly to construct the Markov process associated to the solution X , i.e. a law μ on a space of continuous paths which satisfy the Markov property. Here we see a huge similarity with what we try to do: describe a measure μ which should satisfy a series of constraint (here Markovianity) in an useful and direct way. The solution of the SDE is such a map. Only incidentally it happens in this case that the Brownian motion lives in the same kind of space as the solution. Stochastic calculus is a tool to study the class of stochastic processes arising in this way (and many other related processes arising in more difficult ways). [see e.g. Daniel Revuz and Marc Yor, *Continuous Martingales and Brownian Motion*, 3rd ed. (Springer, 2004).]

Stochastic quantization is here the stochastic analysis of measures of the type encountered in EQFTs.

2 Quantum mechanics

References

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I want to start by giving hints on some aspects of quantum mechanics which are relevant to our discussion. For what concerns us here, quantum mechanics (QM) is the theory which describes certain kind of measurements which cannot be performed with infinite precision and which exhibit a certain complementarity, in the sense that an attempt to a precise measurement of one quantity will result in a larger uncertainty in the measure of another. These strange features concerns primarily with physical phenomena involving interactions at the scales of Planck's quantum of action \hbar . We will however happily ignore the precise value of this constants in what follows.

The basic structure is given by a set \mathcal{O} of observable quantities (e.g. experimental setups, gauges, indicators) and a set \mathcal{S} of states of the physical system. To every observable $A \in \mathcal{O}$ and state $\omega \in \mathcal{S}$ we can associate a real number $\omega(A) \in \mathbb{R}$ which is our representation of the measure of A in the state ω . Operationally this is interpreted in a frequentist way: it corresponds in performing N experiments, each time preparing the system in the fixed state ω and each time measuring the same quantity A in order to obtain the sequence $(x_i)_{i=1, \dots, N}$ of numbers. It is then postulated that the empirical mean converges as $N \rightarrow \infty$ and the limit is $\omega(A)$:

$$\omega(A) = \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{i=1}^N x_i.$$

Given two observables A, B I can measure both and sum the results, obtaining $A + B$ and it is natural to assume that $\omega(A + B) = \omega(A) + \omega(B)$, clearly \mathcal{O} is a (real) vector space and states are linear on \mathcal{O} . Also if I can measure A then I can also measure any (bounded) function $f(A)$ of A by taking $f(x_i)$ each time I observe x_i . It is natural to require that $1 \in \mathcal{O}$ and that $\omega(1) = 1$ for all states ω , moreover states are positive, i.e. $\omega(A) \geq 0$ if the observable A is positive (e.g. if $A = B^2$ for some B). I can measure $A^2, B^2, (A + B)^2$ and define a (commutative) product of two observables $A * B = (A + B)^2 - A^2 - B^2$. The corresponding algebraic structure is called a Jordan algebra. If one requires also that $A_1^2 + \dots + A_n^2 = 0$ implies that $A_1 = \dots = A_n = 0$ then the finite dimensional Jordan algebras can be classified and essentially they are composed of basic blocks given by square self-adjoint matrices with either real, complex, quaternion coefficients, or a finite dimensional Clifford algebra, or the algebra of 3×3 octonionic matrices (the Albert algebra of dimension 27).

Nature seems to have chosen for us only the observable algebras made of complex self-adjoint matrices. In the formalization then one take \mathcal{O} to be the algebra of self-adjoint elements of a C^* -algebra \mathcal{A} , i.e. a Banach algebra endowed with an involution $a \in \mathcal{A} \mapsto a^* \in \mathcal{A}$ which satisfy the condition $\|a^* a\| = \|a\|^2$ (C^* condition). This is true for finite-dimensional complex matrices by taking the involution as Hermitian conjugate and the norm as the operator norm. So seems a good starting point. The theory of C^* algebras is particularly rich and in particular any state on \mathcal{A} is automatically continuous and the norm can be expressed by duality $\|A\| = \sup_{\omega \in \mathcal{S}} |\omega(A)|$.

Any commutative C^* algebra is (up to isomorphism) the algebra of bounded continuous functions on a compact space X (Gelfand's theorem). This implies in particular that for any $x \in X$ we have a state δ_x given by the valuation at the point $x \in X$ (the Dirac measure) such that $\delta_x(AB) = \delta_x(A)\delta_x(B) = A(x)B(x)$. Any other state is a convex combination of such states and we have essentially a probability theory. In particular there exists states which corresponds to infinitely precise measurements of all the observables in $\mathcal{O} \subseteq \mathcal{A}$.

This consequence is at odd with the experience and therefore we have to admit into our modelization non-commutative C^* -algebras. Even in this general case, as soon as we have a sub-algebra which is commutative, it can be essentially considered a system of random variables where any state gives a specific joint probability distribution of the system of observables, in the sense that for any system A_1, \dots, A_n of commutative self-adjoint elements of \mathcal{A} , there exists a measure μ on \mathbb{R}^n such that for all continuous bounded functions $f_1, \dots, f_n: \mathbb{R} \rightarrow \mathbb{R}$ we have

$$\omega(f_1(A_1) \cdots f_n(A_n)) = \int_{\mathbb{R}^n} f_1(x_1) \cdots f_n(x_n) d\mu(x_1, \dots, x_n).$$

From very easy finite-dimensional examples one can construct C^* algebras of observable which are maximally non-commutative, i.e. where there exists states in which one of the observable is precisely determined while another is maximally uncertain (i.e. has uniform distribution). The Gelfand–Naimark theorem states that any C^* algebra can be faithfully represented in Hilbert space \mathcal{H} as a subalgebra of the algebra of all bounded operators $\mathcal{B}(\mathcal{H})$. A more interesting (for us) result is how this representation is obtained, namely via a construction called the Gelfand–Naimark–Segal construction (GNS).

Theorem 4. (GNS) *For any state ω there exists a triplet $(\mathcal{H}_\omega, \pi_\omega, \varphi_\omega)$ where \mathcal{H}_ω is an Hilbert space, π_ω a representation of \mathcal{A} on \mathcal{H}_ω via bounded operators and a vector $\varphi_\omega \in \mathcal{H}_\omega$ such that, for all $A \in \mathcal{A}$*

$$\omega(A) = \langle \varphi_\omega, \pi_\omega(A) \varphi_\omega \rangle.$$

The representation is irreducible iff the state is pure, i.e. if cannot be written as convex linear combination of other states.

A representation π is faithful if $\pi(A) = 0 \Rightarrow A = 0$ and this can be achieved by tensoring all the GNS representations for all the states (it is enough to use the pure states).

Physically the GNS construction is more interesting: it means that the way the quantum system is realized on Hilbert space could apriori depend on the state which we are considering and that different states could result in inequivalent representations, i.e. representations which cannot be unitarily identified one with the other. This is not a particular problem for the QM of finite-dimensional systems since a theorem of von Neumann guarantees that all the faithful representations of the Weyl algebra (the C^* -algebra of a quantum system satisfying certain simple commutation relations) are equivalent, so it is not very important which one we use and only matter of technical convenience. However, as we will see, Poincaré covariance forces us to consider QM with infinitely many degrees of freedom and the question of the inequivalence of representations is at the origin of many problems, in particular of divergences appearing in perturbation theory.

In order to understand better the problem we need to discuss dynamics of a quantum system. So far we essentially only addressed the question of the “kinematics”, i.e. of the description of the motions. For example a classical Newtonian system is completely described by a point in a phase space giving position and momentum of all the particles. The dynamics acts on the phase space as an Hamiltonian flow, in particular it sends the initial positions and momenta (Q_0, P_0) to those at time $t \in \mathbb{R}$: $(Q(Q_0, P_0, t), P(Q_0, P_0, t))$. Any function $f: (Q, P) \mapsto \mathbb{R}$ on the phase space evolves as $f(t, Q_0, P_0) = \alpha_t(f)(Q_0, P_0) = f(Q(Q_0, P_0, t), P(Q_0, P_0, t))$. Note that α_t is an automorphism of the algebra of functions on the phase space which encodes completely the dynamics, moreover $\alpha_{t+s} = \alpha_t \circ \alpha_s$, i.e. $(\alpha_t)_{t \in \mathbb{R}}$ is a continuous one parameter group of automorphism.

In the quantum kinematic picture one has to give away the possibility to describe the motion of the points in the state space (no such points exists since the algebra is non-commutative) and consider the evolution of the observables to be a continuous one-parameter group $(\alpha_t)_{t \in \mathbb{R}}$ acting on the algebra of observables. Let us write $A(t) = \alpha_t(A)$ for the evolution of the observable A . This is called the Heisenberg picture.

The goal of quantum mechanics is therefore to establish which particular triplet $(\mathcal{A}, \omega, \alpha)$ of algebra, state and dynamics allow to predict the behaviour of a given quantum system in a given state. In particular, given the observables A_1, \dots, A_n and the sequence of times t_1, \dots, t_n one would like to compute

$$\omega(\alpha_{t_1}(A_1) \cdots \alpha_{t_n}(A_n)).$$

One of the interests of having a specific representation is that it is then possible to proceed to numerical approximations of such quantities.

The dynamics α is an automorphism, therefore can be made to act on states $(\alpha_t\omega)(A) = \omega(\alpha_t(A))$ since it preserves positivity and normalization, moreover it preserves also purity. It is not however clear if $\alpha_t\omega$ give rise to a GNS representation unitary equivalent to that of ω , i.e. lying on the same folium. Assuming this is the case then one can fix the representation $(\mathcal{H}_\omega, \pi_\omega, \varphi_\omega)$ and associate each state $\alpha_t\omega$ with a unitary transformation $U(t)$ such that $\pi_{\alpha_t\omega}(A) = U(t)^{-1}\pi_\omega(A)U(t)$ and correspondingly $\varphi_{\alpha_t\omega} = U(t)\varphi_\omega$. An important case is when the state is invariant, i.e. $\alpha_t\omega = \omega$ since then $U(t)\varphi_\omega = \varphi_\omega$ for all $t \in \mathbb{R}$ and one can prove that the group $(U(t))_{t \in \mathbb{R}}$ of unitary operators is weakly (and thus strongly) continuous.

The assumption of invariance of the state is important, one can easily construct counterexamples to strong continuity of the group U if it does not hold (see the lecture notes of QMFI, Lecture 17).

Provided we assume the existence of an invariant state ω , then the study of the dynamics α can be reduced to the study of a strongly continuous unitary group U on a fixed Hilbert space $\mathcal{H} = \mathcal{H}_\omega$ with a special vector φ_ω called the vacuum vector or ground state. To each observable a it corresponds a bounded operator $A = \pi_\omega(a)$ and $\pi_\omega(\alpha_t(A)) = U(t)AU(t)^{-1}$, so in particular

$$\omega(\alpha_{t_1}(a_1) \cdots \alpha_{t_n}(a_n)) = \langle \varphi_\omega, A_1 U(t_2 - t_1) A_2 \cdots U(t_n - t_{n-1}) A_n \varphi_\omega \rangle =: W(\{A_k\}, \{t_k\}) \quad (1)$$

where we take $t_n < t_{n-1} < \cdots < t_1$.

We say that φ_ω is cyclic if the span of vectors of the form

$$A_1 U(t_2 - t_1) A_2 \cdots U(t_n - t_{n-1}) A_n \varphi_\omega,$$

for arbitrary A s and times, is dense in \mathcal{H} . Provided the vector φ_ω is cyclic the family of all *correlation functions* of the form (1) for all observables and all increasing sets of times form a complete description of the dynamics of a given quantum system. With our assumptions this system is stationary in time, i.e. the correlation functions depends only on the difference of the set of times.

3 Euclidean quantum mechanics

References

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EQM means to do QM in “imaginary time”. What does this means precisely? We need to understand what is the unitary group $U(t)$ at imaginary times.

By Stone's theorem the strongly continuous unitary group $(U(t))_{t \in \mathbb{R}}$ corresponds to a homomorphism $X: C(\mathbb{R}; \mathbb{C}) \rightarrow \mathcal{B}(\mathcal{H})$ of C^* -algebras such that $X(e^{it\cdot}) = U(t)$ for all $t \in \mathbb{R}$.

We say that U is of positive energy if $X(f) = 0$ for all f with support in $\{x < 0; x \in \mathbb{R}\}$. In this case one can define $K(s) = X(e^{-s\cdot} \mathbb{1}_{\mathbb{R}_+})$ for all $s \geq 0$ (with some care) and observe that $K(t+s) = K(t)K(s)$, $\|K(t)\| \leq 1$ and that $t \mapsto K(t)$ is strongly continuous. So $(K(t))_{t \geq 0}$ is a strongly continuous semigroup of contractions.

The notable fact is that, given $(K(t))_{t \geq 0}$ one can reconstruct X and then U so they correspond to each other one-to-one and express essentially the same object, in our case the Hilbert space realisation of the dynamics of the quantum system.

The idea now is to take correlation functions for a dynamics with positive energy and continue it to imaginary time differences to obtain correlation functions of the form

$$S(\{A_k\}, \{t_k\}) := \langle \varphi_\omega, A_1 K(t_1) A_2 \cdots K(t_{n-1}) A_n \varphi_\omega \rangle \quad (2)$$

for arbitrary operators A_k and positive times t_k . The question is now whether we can recover the full structure of the quantum problem only knowing these functions. In particular we would like to reconstruct the Hilbert space (with all its vectors) and the scalar product, and then also the unitary dynamics U .

We need a way to decode from these functions the key properties which allow the reconstruction. Surely we need the fact that K is a contractive semigroup. A more subtle property comes from the positivity of the Hilbert scalar product, indeed let us observe that we always have, for example that the correlation function

$$\begin{aligned} S((A_3, A_2, A_1, A_1, A_2, A_3), (t_2, t_1, t_1, t_2)) &:= \langle \varphi_\omega, A_3 K(t_2) A_2 K(t_1) A_1 A_1 K(t_1) A_2 K(t_2) A_3 \varphi_\omega \rangle \\ &= \langle A_1 K(t_1) A_2 K(t_2) A_3 \varphi_\omega, A_1 K(t_1) A_2 K(t_2) A_3 \varphi_\omega \rangle \geq 0 \end{aligned}$$

is positive. Of course similar considerations hold for similar functions and for certain linear combinations of such functions (guess their form). It could be formulated explicitly but the formula will not give any further insight, so we avoid it here. This property is called reflection positivity of the family of correlation functions 2. Let us give also a name to such a family $\{S(\mathbb{A}, \mathbb{T})\}_{\mathbb{A}, \mathbb{T}}$ and call it Schwinger functions. Here we denote \mathbb{A} an arbitrary vector of observables and \mathbb{T} one of times. Of course here it is trivial that such a property holds, but the idea is that we will call Schwinger functions all systems of functions indexed by \mathbb{A}, \mathbb{T} and satisfying reflection positivity and some other conditions encoding the semigroup property of K (compatibility condition) and an analytic condition which encodes the contractivity of the semigroup K and which essentially says that Schwinger functions have estimates

$$\left| \int S((A_k), (t_1, \dots, t_n)) g_1(t_1) \cdots g_n(t_n) dt_1 \cdots dt_n \right| \leq \left(\prod_k \|A_k\| \right) \prod_k \left(\sup_{x \in \mathbb{R}_+} \left| \int e^{-tx} g_k(t) dt \right| \right). \quad (3)$$

This condition implies that $(t_1, \dots, t_n) \mapsto S((A_k), (t_1, \dots, t_n))$ is the (multidimensional) Laplace transform of a tempered distribution (see e.g. B. Simon “The $P(\varphi)_2$ Euclidean Quantum Field Theory”, Chap. 2, Sect. 2.2). (This is an unnecessarily strong bound and a weaker condition would be sufficient).

The baby version of the Osterwalder–Schrader reconstruction theorem says that from the family of Schwinger functions $\{S(\mathbb{A}, \mathbb{T})\}_{\mathbb{A}, \mathbb{T}}$ one can recover the QM data, in particular the Wightman functions $\{W(\mathbb{A}, \mathbb{T})\}_{\mathbb{A}, \mathbb{T}}$. The idea of the proof is to use a GNS construction on the Schwinger functions to obtain the Hilbert space \mathcal{H} , the vacuum, the representation π of \mathcal{A} and a semigroup K . In order to prove the continuity and contractivity of the semigroup the analytic condition (3) is needed. From there one can then construct the positive energy group U and therefore the original QM dynamical problem.

So the problem of the construction of interesting QM dynamics is reduced to that of finding Schwinger functions, i.e. correlation functions which satisfy the three above mentioned conditions.

It is a remarkable fact that the correlation functions of certain probabilistic models are Schwinger functions. To show this in a simple setting we assume that the observables indexing the Schwinger functions belong to a commutative subalgebra \mathcal{O}' isomorphic to $C(\mathbb{R})$ (i.e. we have only one continuous degree of freedom) so that we can identify an element A with the function $a \in C(\mathbb{R})$. Then we can associate to a real valued stochastic process $(X_t)_{t \in \mathbb{R}}$ the family of Schwinger functions

$$S(a_1, \dots, a_n, t_1, \dots, t_{n-1}) = \mathbb{E}[a_1(X_{s_1}) \cdots a_{n-1}(X_{s_{n-1}}) a_n(X_T)]$$

with $s_k = t_k + \dots + t_{n-1} + T$ for an arbitrary $T \in \mathbb{R}$. In order for this definition not to depend on T we require that the process X is stationary in time, i.e. the processes $(X_t)_{t \in \mathbb{R}}$ and $(X_{t+s})_{t \in \mathbb{R}}$ have the same law for all $s \in \mathbb{R}$. So the question now becomes: under which conditions on the law of X these Schwinger functions have the form (2) for some quantum data and contractive semigroup K ?

It is not difficult to show that the family of such functions satisfy the compatibility conditions required by the structure of the (so far only putative) r.h.s. of (2).

Reflection positivity is trickier: the explicit form of the functions S give us a simpler way to write the RP condition. On functions F on $C(\mathbb{R}; \mathbb{R})$ we can introduce an operation Θ of time inversion such that $(\Theta F)(x) = F(\theta x)$ with $(\theta x)(t) = x(-t)$. Then for any complex function F on $C(\mathbb{R}_{\geq 0}; \mathbb{R})$ of the form

$$F(x) = \sum_k c_k e^{i\lambda_k x(t_k)} \quad (4)$$

with coefficients $c_k \in \mathbb{C}$, $\lambda_k \in \mathbb{R}$, $t_k \in \mathbb{R}_{\geq 0}$ (note that the times are positive!), the RP of the Schwinger functions become the relation

$$\mathbb{E}[\overline{(\Theta F)(X)} F(X)] \geq 0. \quad (5)$$

Definition 5. A measure on $C(\mathbb{R}; \mathbb{R})$ is reflection positive iff eq. (5) holds for any cylinder function.

This is already a nontrivial condition. For example, if X is a Gaussian process we can use cylinder functions to approximate any bounded continuous function and then also polynomials and therefore obtain that a reflection positive Gaussian measure should be centred and with covariance $C(t) = \mathbb{E}[X_t X_0]$ satisfying

$$0 \leq \mathbb{E} \left[\left(\sum_k \bar{c}_k X_{-t_k} \right) \left(\sum_k c_k X_{t_k} \right) \right] = \sum_{k,k'} c_k \bar{c}_{k'} C(t_k + t_{k'})$$

for all choices of coefficients here. Functions satisfying this property are called totally monotone. By a theorem of Bernstein a bounded and totally monotone function $C: \mathbb{R}_{\geq 0} \rightarrow \mathbb{C}$ has the representation

$$C(t) = \int_{\mathbb{R}_{\geq 0}} e^{-tx} \nu(dx)$$

for some bounded positive measure ν on \mathbb{R}_+ . In this case, indeed,

$$\sum_{k,k'} c_k \bar{c}_{k'} C(t_k + t_{k'}) = \int_{\mathbb{R}_{\geq 0}} \left| \sum_k c_k e^{-t_k x} \right|^2 \nu(dx) \geq 0.$$

It turns out that this property is sufficient to have reflection positivity of the Gaussian process X . The simplest case is when the measure ν is concentrated in a point $\alpha > 0$, then we have the covariance

$$C(t) = \frac{e^{-\alpha|t|}}{2\alpha}, \quad t \in \mathbb{R},$$

with an arbitrary normalization (see below). The corresponding Gaussian process is called the Ornstein–Uhlenbeck process (OU) and we conclude that any scalar reflection positive Gaussian process in one dimension can be constructed by taking sums of independent OU processes (see the QMFI lecture notes).

3.1 Symmetric Markov processes

Another important strategy to obtain RP processes is to use Markovianity. Let X be a Markov process, i.e. such that

$$\mathbb{E}[f(X_t) | \mathcal{F}_s] = \mathbb{E}[f(X_t) | X_s]$$

where $(\mathcal{F}_s = \sigma(X_r; r \leq s))_{s \in \mathbb{R}}$ is the filtration generated by X . Recalling the support of the function F in the RP condition, we have

$$\mathbb{E}[\overline{\Theta F(X)} F(X)] = \mathbb{E}[\overline{\Theta F(X)} \mathbb{E}[F(X)|\mathcal{F}_0]] = \mathbb{E}[\overline{\Theta F(X)} \mathbb{E}[F(X)|X_0]] = \mathbb{E}[\overline{\mathbb{E}[\Theta F(X)|X_0]} \mathbb{E}[F(X)|X_0]]$$

an easy way to force this quantity to be positive is then to take $\mathbb{E}[\Theta F|X_0] = \mathbb{E}[F|X_0]$ which is easily seen to imply that the law of X is invariant under the time reversal Θ . Stationary and time-reversal invariant Markov processes are then reflection positive. The converse is also true: RP Markov processes are time-reversal invariant and stationary.

Having a symmetric (i.e. time-reversal invariant) Markov process allows to represent the reconstruction procedure quite explicitly. The idea is to consider the complex vector space \mathcal{E}_+ of cylindrical function supported on positive times as introduced above and make a pre-Hilbert space out of it by considering the Hermitian scalar product

$$\langle F, G \rangle = \mathbb{E}[\overline{\Theta F} G]$$

which is positive definite thanks to the RP of the law of the process X . We take the quotient $\mathcal{E}_+ \setminus \mathcal{N}$ where \mathcal{N} is the subspace of elements in \mathcal{E}_+ with zero norm and complete wrt. the scalar product to obtain a Hilbert space \mathcal{H} . On this Hilbert space the elements $a \in \mathcal{A}$ act as multiplication operators via $(af)(X) = a(X_0)f(X)$ which is a C^* -homomorphism. Moreover we have the natural action of time-translation via $t \geq 0$: T_t which acts on \mathcal{E}_+ (since for example $T_t X_s = X_{s+t}$ with $s+t \geq 0$ if $s, t \geq 0$). Moreover we have that elements in the form $F - \mathbb{E}[F|X_0]$ have zero norm, since by Markovianity

$$\langle F, \mathbb{E}[F|X_0] \rangle = \mathbb{E}[\overline{\Theta F} \mathbb{E}[F|X_0]] = \mathbb{E}[\overline{\Theta F} \mathbb{E}[F|\mathcal{F}_0]] = \langle F, F \rangle$$

so $\|F - \mathbb{E}[F|X_0]\| = 0$. By stationarity the operators T_t form a symmetric semigroup:

$$\langle T_t F, G \rangle = \mathbb{E}[\overline{\Theta T_t F} G] = \mathbb{E}[T_{-t} \overline{\Theta F} G] = \mathbb{E}[\overline{\Theta F} T_t G] = \langle F, T_t G \rangle$$

which preserves the null space \mathcal{N} since if $F \approx 0$ then by this formula we see that $\langle T_t F, G \rangle = 0$ for all G , i.e. $T_t F \approx 0$. It is also a contractive operator for all $t \geq 0$:

$$\langle T_t F, T_t F \rangle = \langle F, T_{2t} F \rangle \leq \|F\|^{1/2} \|T_{2t} F\|^{1/2}$$

and by iterating this inequality we get

$$\|T_t F\|^2 \leq \|F\|^{1/2+1/2^2} \|T_{2^2 t} F\|^{1/2^2} \leq \dots \leq \|F\|^{1/2+\dots+1/2^n} \|T_{2^n t} F\|^{1/2^n} \leq \|F\|$$

since

$$\|T_{2^n t} F\|^2 = \mathbb{E}[T_{-2^n t} \overline{\Theta F} T_{2^n t} F] \leq (\mathbb{E}[|T_{-2^n t} \overline{\Theta F}|^2])^{1/2} (\mathbb{E}[|T_{2^n t} F|^2])^{1/2} = (\mathbb{E}[|\overline{\Theta F}|^2])^{1/2} (\mathbb{E}[|F|^2])^{1/2}$$

which is independent of n .

Strong continuity of $(T_t)_t$ follows then via approximations from weak continuity and from the fact that the process X is continuous in distribution (Exercise).

We have already shown that $F - \mathbb{E}[F|X_0] \in \mathcal{N}$ so in the completed, quotiented Hilbert space \mathcal{H} we have that all the vectors can be represented as $\mathbb{E}[F|X_0]$ for some $F \in L^2(\mathbb{P})$ measurable wrt $\sigma(X_s; s \geq 0)$. Therefore really all these vectors are in $L^2(\mathbb{R}, \rho)$ where ρ is the law of X_0 :

$$\begin{aligned} \langle F, G \rangle_{\mathcal{H}} &= \mathbb{E}[\overline{\Theta F} G] = \mathbb{E}[\overline{\Theta F} \mathbb{E}[G|\mathcal{F}_0]] = \mathbb{E}[\overline{\mathbb{E}[F|X_0]} \mathbb{E}[G|X_0]] \\ &= \mathbb{E}[\overline{f(X_0)} g(X_0)] = \int_{\mathbb{R}} \overline{f(x)} g(x) \rho(dx) = \langle f, g \rangle_{L^2(\mathbb{R}, \rho)} \end{aligned}$$

with $f, g: \mathbb{R} \rightarrow \mathbb{C}$ such that $f(X_0) = \mathbb{E}[f|X_0]$ and $g(X_0) = \mathbb{E}[g|X_0]$ \mathbb{P} -almost surely. This shows that \mathcal{H} can be isometrically embedded in $L^2(\mathbb{R}, \rho)$, the converse is also true. In this representation the semigroup $(T_t)_{t \geq 0}$ give rise to a semigroup $(K(t))_{t \geq 0}$ on $L^2(\mathbb{R}, \rho)$, indeed for any $f \in L^2(\mathbb{R}, \rho)$ we can consider $T_t(f(X_0)) = f(X_t)$ and then condition again onto X_0 and define

$$(K(t)f) = h, \quad h(X_0) = \mathbb{E}[f(X_t)|X_0].$$

so that $\langle F, T_t G \rangle_{\mathcal{H}} = \langle f, K(t)g \rangle_{L^2(\mathbb{R}, \rho)}$.

From this contractive semigroup one then construct a unitary dynamics U and therefore a model of QM dynamics. These considerations can also be extended without essential changes to the case of infinite-dimensional Markov processes, generalisation needed in QFT. It has been Nelson which first suggested the use of Markov processes to construct models of QFT along the lines we described. Unfortunately the infinity of degrees of freedom makes difficult to prove Markovianity of interesting examples. And it was realised by Osterwalder and Schröder that the weaker property of RP is much easier to prove and enough for the reconstruction of the QM data.

3.2 An important example and some Gaussian analysis

I think now useful to specialize the above considerations to an important example given by the OU process already introduced before which is RP positive.

Notably the OU process is Markovian, indeed is the solution to a linear SDE involving a Brownian motion B :

$$dX_t = -\alpha X_t dt + dB_t, \quad X_0 \sim \mathcal{N}(0, (2\alpha)^{-1})$$

indeed the explicit solution of this equation, via variation of constants is

$$X_t = e^{-\alpha(t+S)} X_{-S} + \int_{-S}^t e^{-\alpha(t-s)} dB_s,$$

where the stochastic integral is a Wiener integral (i.e. the limit in $L^2(\mathbb{P})$ of Riemann sums). Taking $S \rightarrow \infty$ the first term disappear and one remains with the Gaussian process

$$X_t = \int_{-\infty}^t e^{-\alpha(t-s)} dB_s$$

which can be checked to have the correct covariance, to be stationary and to have invariant measure $\mathcal{N}(0, (2\alpha)^{-1})$. By Ito's formula, for $t \geq t'$

$$\mathbb{E}[X_t X_{t'}] = \mathbb{E} \left[\left(\int_{-\infty}^t e^{-\alpha(t-s)} dB_s \right) \left(\int_{-\infty}^{t'} e^{-\alpha(t'-s)} dB_s \right) \right] = e^{-\alpha(t-t')} \int_{-\infty}^{t'} e^{-2\alpha(t'-s)} ds = \frac{e^{-\alpha(t-t')}}{2\alpha}.$$

Note that the time-reversal symmetry is not quite explicit in the SDE formulation since Ito calculus breaks the invariance.

From the SDE we have that

$$X_t = e^{-\alpha t} X_0 + \int_0^t e^{-\alpha(t-s)} dB_s = e^{-\alpha t} X_0 + \mathcal{N} \left(0, \frac{1 - e^{-2\alpha t}}{2\alpha} \right)$$

with the Gaussian r.v. independent of \mathcal{F}_0 , so in particular it is a Markov process and

$$(K(t)f)(x) = \mathbb{E}[f(X_t)|X_0 = x] = \mathbb{E} \left[f \left(e^{-\alpha t} x + \mathcal{N} \left(0, \frac{1 - e^{-2\alpha t}}{2\alpha} \right) \right) \right] = \mathbb{E} \left[f \left(e^{-\alpha t} x + \left(\frac{1 - e^{-2\alpha t}}{2\alpha} \right)^{1/2} Z \right) \right]$$

where $Z \sim \mathcal{N}(0, 1)$. One can check explicitly that this is a semigroup on $L^2(\mathbb{R}, \rho)$ which is self-adjoint, positive (i.e. $f \geq 0 \Rightarrow K(t)f \geq 0$), contractive, strongly-continuous. It is an important ingredient in the analysis of Gaussian r.v. (both in finite and infinite dimensions). One has also

$$K(t)f \rightarrow \int f d\rho$$

in $L^2(\rho)$ as $t \rightarrow \infty$ (the convergence holds also in stronger senses).

If we want to understand better the associated quantum dynamics in $L^2(\rho)$, we need a grasp of the unitary operator U associated to K via the procedure we described above. We could compute the generator L of K and prove that it is a positive self-adjoint operator and then consider the associated unitary group.

Another way to proceed in this very explicit situation is just to diagonalize K . Which is what we are going to do now. Let $e_\lambda(x) = \exp(i\lambda x)$ the trigonometric exponential, then

$$\begin{aligned} (K(t)e_\lambda)(x) &= \mathbb{E} \left[e_\lambda \left(e^{-at}x + \left(\frac{1-e^{-2at}}{2\alpha} \right)^{1/2} Z \right) \right] = e_\lambda(e^{-at}x) \mathbb{E} \left[e_\lambda \left(\left(\frac{1-e^{-2at}}{2\alpha} \right)^{1/2} Z \right) \right] \\ &= e_\lambda(e^{-at}x) \exp \left(-\frac{\lambda^2}{2} \left(\frac{1-e^{-2at}}{2\alpha} \right) \right) \end{aligned}$$

so if we let $\hat{e}_\lambda(x) = \exp(i\lambda x) \exp \left(\frac{\lambda^2}{2} \frac{1}{2\alpha} \right)$ we have

$$(K(t)\hat{e}_\lambda)(x) = \hat{e}_\lambda e^{-at}(x)$$

By expanding both the l.h.s. and r.h.s in powers of λ we obtain

$$\sum_{n \geq 0} \frac{(i\lambda)^n}{n!} K(t)H_n(x) = (K(t)\hat{e}_\lambda)(x) = \hat{e}_\lambda e^{-at}(x) = \sum_{n \geq 0} \frac{(i\lambda)^n}{n!} e^{-ant} H_n(x)$$

where we denoted $H_n(x)$ the coefficients in this expansion. Each H_n is a polynomial in x with α dependent coefficients and highest order term x^n , they are called Hermite polynomials. We conclude that

$$K(t)H_n(x) = e^{-ant} H_n(x), \quad t \geq 0,$$

i.e. Hermite polynomials are eigenfunctions of each $K(t)$. Since $K(t)$ is symmetric this proves also that the $(H_n)_n$ are orthogonal (recall that we are here in $L^2(\rho)$) since they corresponds to different eigenvalues of $K(t)$. It is not difficult also to show that they are total in $L^2(\rho)$, i.e. every vector can be written as a convergent sum of Hermite polynomials, indeed assume that this is not the case for the non-null vector f , then $\langle \hat{e}_\lambda, f \rangle = 0$ for all $\lambda \in \mathbb{R}$, but this implies that $\langle e_\lambda, f \rangle = 0$ since the two functions differ by a strictly positive constant factor. The functions e_λ can approximate on any compact set any arbitrary continuous function in the uniform norm, so by an approximation procedure we can show that $\langle g, f \rangle = 0$ for all $g \in L^2(\rho)$ and we conclude that $f = 0$ (we could have used also the Fourier transform to conclude).

At this point it is easy to construct the unitary group $U(t)$ (since we can reason in a one dimensional Hilbert space), we have:

$$U(t)H_n = e^{iant} H_n, \quad t \in \mathbb{R}.$$

We have completely solved the QM problem for this model: we have an Hilbert space $L^2(\rho)$, a representation of $C(\mathbb{R})$ given by multiplication operators, a strongly continuous unitary dynamical group $(U(t))_{t \in \mathbb{R}}$ and a ground state given by the constant function 1.

From U (or K) we can produce the homeomorphism $C(\mathbb{R}_+) \rightarrow \mathcal{B}(L^2(\rho))$ explicitly

$$E(f)H_n = f(\alpha n)H_n$$

and one see that it is actually an homeomorphism from $C(\alpha\mathbb{N})$. This new commutative subalgebra of $\mathcal{B}(L^2(\rho))$ does not commute with the algebra $\mathcal{Q} \approx C(\mathbb{R})$ (which acts as multipliers), moreover it corresponds to an observable which can take only discrete integer values. The system is quantized. The observables $X(f)$ do not change in time, by construction $U(t)^{-1}X(f)U(t) = X(f)$: they are constants of motion (it is indeed the “energy” of the system).

Note that $\partial_x \hat{e}_\lambda(x) = i\lambda \hat{e}_\lambda(x)$ which means that $\partial_x H_n = nH_{n-1}$, on the other hand by an integration by parts:

$$\delta_{m,n} \int H_m H_n d\rho = \frac{1}{n+1} \int H_m \partial_x H_{n+1} d\rho = \frac{1}{n+1} \int H_{n+1} \left(-\frac{x}{2\alpha} - \partial_x \right) H_m d\rho$$

so

$$\left(\frac{x}{2\alpha} - \partial_x \right) H_n = c_{n+1} H_{n+1} \quad \int H_n H_n d\rho = \frac{c_{n+1}}{n+1} \int H_{n+1} H_{n+1} d\rho$$

to determine c_{n+1} we compute

$$c_{n+1} H_n = \frac{c_{n+1}}{n+1} \partial_x H_{n+1} = \frac{1}{n+1} \partial_x \left(\frac{x}{2\alpha} - \partial_x \right) H_n = \frac{1}{n+1} \left[\frac{1}{2\alpha} H_n + \left(\frac{x}{2\alpha} - \partial_x \right) \partial_x H_n \right] = \frac{1}{n+1} \left[\frac{1}{2\alpha} H_n + n c_n H_n \right]$$

which gives $n c_n = n/2\alpha$ since $c_1 = 1/2\alpha$. Let us register them here by giving names to these operators

$$A = \partial_x, \quad C = \left(\frac{x}{2\alpha} - \partial_x \right)$$

are called annihilation and creation operators. (They are unbounded but with a very simple algebraic structure over the Hermite basis so not so difficult to analyze. We will not use them much here, apart from specific algebraic computations on the Hermite basis). Note that $C^* \supseteq A$. We discovered several relations for them

$$A H_n = n H_{n-1}, \quad C H_n = \frac{1}{2\alpha} H_{n+1},$$

$$C A H_n = \left(\frac{x}{2\alpha} - \partial_x \right) \partial_x H_n = \frac{n}{2\alpha} H_n.$$

By putting them together we have also

$$\frac{x}{2\alpha} H_n = n H_{n-1} + \frac{1}{2\alpha} H_{n+1}$$

and in particular that

$$\left(\frac{x}{2\alpha} - \partial_x \right)^n 1 = \frac{1}{(2\alpha)^n} H_n.$$

Moreover

$$[A, C] = \left[\partial_x, \left(\frac{x}{2\alpha} - \partial_x \right) \right] = \partial_x \left(\frac{x}{2\alpha} - \partial_x \right) - \left(\frac{x}{2\alpha} - \partial_x \right) \partial_x = \frac{1}{2\alpha}$$

We have also $x = 2\alpha(C + A)$. Let's then define

$$P = (2\alpha i)(C - A) = i(x - 2(2\alpha)\partial_x), \quad Q = 2\alpha(C + A),$$

and observe that P is symmetric and that

$$Q^2 = (2\alpha)^2 (C + A)^2 = (2\alpha)^2 (C^2 + CA + AC + A^2)$$

$$P^2 = -(2\alpha)^2 (C - A)^2 = -(2\alpha)^2 (C^2 + A^2 - AC - CA)$$

so

$$P^2 + Q^2 = 2(2\alpha)^2(CA + AC) = 2(2\alpha^2)(CA + AC) = 2(2\alpha)^2\left(2CA + \frac{1}{2\alpha}\right) = 2(2\alpha) + 4(2\alpha)^2CA$$

From which one gets

$$\left[\frac{1}{8}(P^2 + Q^2) - \frac{1}{2}\alpha\right]H_n = (\alpha n)H_n$$

so we can identify the generator $E \geq 0$ of the dynamics of system with

$$E = \frac{1}{8}(P^2 + Q^2) - \frac{1}{2}\alpha \geq 0$$

which resembles the classical Hamiltonian for the Harmonic oscillator, note that

$$[P, Q] = [(2\alpha i)(C - A), 2\alpha(C + A)] = -i4\alpha$$

so α describes the degree of non-commutativity of the algebra.

The us wrap up a take-away message from all these computations. A first one is that even if we restricted our considerations to a particular commutative subalgebra Q of \mathcal{A} the fact that we included the dynamics in the description carry on information on the non-commutative nature of our original problem in the time evolution. In particular we can construct other subalgebras (those of the bounded functions of P or E) which do not commute with the Q algebra. Note also that the algebra generated by P, Q contains that of E so also the dynamics and can be considered therefore a complete description of the quantum system.

This dynamics is really boring however. Every pure state ω (in the folium of the stationary state ω_0) of the system corresponds to a vector f of the Hilbert space (any other state can be written as a linear combination of pure states). Every vector f can be decomposed on the basis (H_n) on which the dynamics is linear, so for an observable $a \in C(\mathbb{R})$ we have

$$\begin{aligned} \omega_0(\alpha_t(a) * a) &= 2\text{Re}\langle 1, U(t)^{-1}Q(a)U(t)Q(a)1 \rangle \\ &= 2\text{Re}\sum_n e^{i\alpha t n} \langle 1, Q(a)H_n \rangle \langle H_n, Q(a)1 \rangle = 2\text{Re}\sum_n e^{i\alpha t n} |\langle H_n, Q(a)1 \rangle|^2, \end{aligned}$$

which could be interpreted by saying that the first measurement of a produces n particles from the vacuum state which then evolve for time t via simple oscillations and are reabsorbed by the second measurement with exactly the same amplitude.

By an approximation procedure, one can polarize the above expression with two observables a, b and by an approximation procedure take $Q(a) \approx H_n$ and $Q(b) \approx H_m$, in which case one gets

$$\omega_0(\alpha_t(a) * b) \approx 2\text{Re}\sum_k e^{i\alpha t k} \langle 1, Q(a)H_k \rangle \langle H_k, Q(b)1 \rangle = \delta_{n=m} 2\text{Re}(e^{i\alpha t n}) |\langle H_n, H_n \rangle|^2$$

which can be interpreted by saying that if one creates n particles at time 0 and let them evolve up to time t then they remains n particles and one will never measure m particles if $m \neq n$.

Another important take-away for us is that these computations give us tools for the Gaussian analysis later on.

3.3 Perturbations of RP processes

References

- József Lőrinczi, Fumio Hiroshima, and Volker Betz, *Feynman-Kac-Type Theorems and Gibbs Measures on Path Space: With Applications to Rigorous Quantum Field Theory*, De Gruyter Studies in Mathematics 34 (Berlin; Boston: De Gruyter, 2011).

As we have said, it is useful to have a source of RP processes which do not rely on Markovianity and survive the generalisation to infinite dimensions which we are going to pursue later on.

A convenient way to construct a large class of RP processes is to take Gibbsian perturbations of a RP process. With this we mean consider a potential function $V: \mathbb{R} \rightarrow \mathbb{R}_{\geq 0}$ and a new probability measure \mathbb{Q} given by

$$\mathbb{Q}_T(d\omega) = \frac{1}{Z_T} \exp\left(-\int_{-T}^T V(X_s(\omega)) ds\right) \mathbb{P}(d\omega)$$

with a normalization factor Z_T .

Lemma 6. *The measure \mathbb{Q}_T is reflection positive.*

Proof. Take

$$G = \exp\left(-\int_0^T V(X_s) ds\right),$$

then

$$\Theta G = \exp\left(-\int_0^T V(X_{-s}) ds\right) = \exp\left(-\int_{-T}^0 V(X_s) ds\right)$$

and

$$\exp\left(-\int_{-T}^T V(X_s(\omega)) ds\right) = (\Theta G)G = (\overline{\Theta G})G$$

then for any $F \in \mathcal{E}_+$ we have

$$\begin{aligned} \mathbb{E}_{\mathbb{Q}_T}[\overline{\Theta F} F] &= \frac{1}{Z_T} \mathbb{E}_{\mathbb{P}}\left[\overline{\Theta F} F \exp\left(-\int_{-T}^T V(X_s(\omega)) ds\right)\right] \\ &= \frac{1}{Z_T} \mathbb{E}_{\mathbb{P}}[\overline{\Theta F} F (\overline{\Theta G})G] = \frac{1}{Z_T} \mathbb{E}_{\mathbb{P}}[\overline{\Theta(FG)}(FG)] \geq 0 \end{aligned}$$

since \mathbb{P} is RP. □

Remark. More “fancy” perturbations, like e.g.

$$\exp\left(-\int_{-T}^T \int_{-T}^T W(X_s(\omega), X_{s'}(\omega)) ds ds'\right),$$

are not in general reflection positive. Note also that to prove RP we used essentially that there is an integral over time and the multiplicativity of the exponential.

The problem is that the measure \mathbb{Q}_T is not stationary anymore and to obtain a stationary measure we would need to take the limit $T \rightarrow \infty$, indeed if T_t denote the time translation on functions on the path space $T_t F(x) = F(T_t x)$ with $(T_t x)(s) = x(t+s)$, then by stationarity of \mathbb{P} we have

$$\begin{aligned} \mathbb{E}_{\mathbb{Q}_T}[T_t F] &= \frac{1}{Z_T} \mathbb{E}_{\mathbb{P}}\left[T_t F \exp\left(-\int_{-T}^T V(X_s(\omega)) ds\right)\right] = \frac{1}{Z_T} \mathbb{E}_{\mathbb{P}}\left[T_t F T_t \exp\left(-\int_{-T}^T V(X_{s-t}(\omega)) ds\right)\right] \\ &= \frac{1}{Z_T} \mathbb{E}_{\mathbb{P}}\left[F \exp\left(-\int_{-T}^T V(X_{s-t}(\omega)) ds\right)\right] = \frac{1}{Z_T} \mathbb{E}_{\mathbb{P}}\left[F \exp\left(-\int_{-T-t}^{T-t} V(X_s(\omega)) ds\right)\right] \end{aligned}$$

which is not what we want and we would need to take the limit $T \rightarrow \infty$ in both sides to obtain that $\mathbb{E}_{\mathbb{Q}_\infty}[T_t F] = \mathbb{E}_{\mathbb{Q}_\infty}[F]$. The existence and uniqueness of such limits are nontrivial and depends on the specific form of the potential V (and on the law of X under \mathbb{P} of course).

Let us make connection with the Markovian point of view. Let X be a RP Markov process with semigroup K and introduce a new semigroup Q^V (not necessarily contractive) via the Feynman–Kac–Nelson formula

$$(Q_t^V f)(x) = \mathbb{E}_{\mathbb{P}} \left[f(X_t) \exp \left(- \int_0^t V(X_s) ds \right) \middle| X_0 = x \right],$$

where we assume we have a regular version of the conditional probability $\mathbb{P}(*|X_0)$ (this requires some mild assumptions on X of course). That it is a semigroup derive easily from the Markovianity of X (exercise). It is positive and in particular satisfies $|(Q_t^V f)(x)| \leq (Q_t^V |f|)(x)$, moreover it is symmetric

$$\langle g, Q_t^V f \rangle = \mathbb{E}_{\mathbb{P}} \left[\overline{g(X_0)} \exp \left(- \int_0^t V(X_s) ds \right) f(X_t) \right] = \mathbb{E}_{\mathbb{P}} \left[\overline{g(X_{-t})} \exp \left(- \int_{-t}^0 V(X_s) ds \right) f(X_0) \right] = \langle Q_t^V g, f \rangle.$$

Assume that $\|Q_s^V\| < \infty$ for some $s > 0$. Then

$$|\langle g, Q_{2s}^V f \rangle| = |\langle Q_s^V g, Q_s^V f \rangle| \leq \|Q_s^V\|^2 \|g\| \|f\| \Rightarrow \|Q_{2s}^V\| \leq \|Q_s^V\|^2$$

and conversely

$$|\langle Q_{s/2}^V g, Q_{s/2}^V f \rangle| = |\langle g, Q_s^V f \rangle| \leq \|Q_s^V\| \|g\| \|f\| \Rightarrow \|Q_{s/2}^V\|^2 \leq \|Q_s^V\|$$

so iterating these inequalities we have $\|Q_{2^n s}^V\| \leq \|Q_s^V\|^{2^n}$, and $\|Q_{s/2^n}^V\| \leq \|Q_s^V\|^{2^{-n}}$. For any $t \geq 0$ we can write down the dyadic decomposition of $p = t/s$ and get $\|Q_t^V\| \leq \|Q_s^V\|^p \leq \|Q_s^V\|^{t/s}$. We have proven that there exists a constant $c = \|Q_s^V\|^{1/s}$ such that $\|Q_t^V\| \leq e^{ct}$. By repeating the argument with t we have now $\|Q_s^V\|^{1/s} \leq \|Q_t^V\|^{1/t} \leq \|Q_s^V\|^{1/s}$ so we must have $\|Q_t^V\| = e^{ct}$.

We can then define $K^V(t)f := e^{-ct} Q_t^V f$ and obtain a contractive semigroup $K^V(t)$ for which $\|K^V(t)\| = 1$ for all $t \geq 0$, in particular there must exist at least one vector ψ such that $K^V(t)\psi = \psi$ and we are back in the situation where we have a ground state ψ and a contractive semigroup. If we assume that this eigenvector is unique (see e.g. the Perron–Frobenius theorem) then one can prove that for any other vector f we have

$$K^V(t)f \rightarrow \langle \psi, f \rangle \psi$$

and in particular $K^V(t)1 \rightarrow \langle \psi, 1 \rangle \psi$. Note that we can assume ψ positive since $|\psi| = |K^V(t)\psi| \leq K^V(t)|\psi|$

$$\langle |\psi|, K^V(t)|\psi| \rangle \geq \langle |\psi|, |K^V(t)\psi| \rangle \geq \langle \psi, K^V(t)\psi \rangle = \langle \psi, \psi \rangle = 1$$

so $|\psi|$ is also an eigenvector with eigenvalue 1. (We have also to prove that does not depend on t , exercise).

One situation when one has a unique ground state is when the operator K^V is positivity improving, i.e. when for some $t > 0$ and for any $f \geq 0$ such that $f \neq 0$ we have $K^V(t)f(x) > 0$ for all x .

One can prove that in the situation where there is uniqueness of the groundstate the Gibbsian recipe above give exactly as a limit the measure \mathbb{Q} corresponding to the stationary Markov process with transition operator given by K^V .

3.4 Euclidean Quantum Field Theory

References

- R. Haag, *Local Quantum Physics: Fields, Particles, Algebras*, 2nd rev. and enl. ed, Texts and Monographs in Physics (Springer, 1996), <http://gen.lib.rus.ec/book/index.php?md5=E91268014A4AF250E095387BD5C2A678>.
- F. Strocchi, *An Introduction to Non-Perturbative Foundations of Quantum Field Theory* (Oxford: OUP Oxford, 2013).

How we extend all these considerations to QFTs instead of just QM models? This will involve to consider QM with infinitely many degrees of freedom. (Recall that so far we gave only explicit examples with one degree of freedom, i.e. we were measuring only one quantity at any given time $\text{Hom}(C(\mathbb{R}, \mathbb{C}), \mathcal{A})$).

In $(n+1)$ -dimensional Minkowski space \mathbb{M}^{n+1} the situation looks like this:

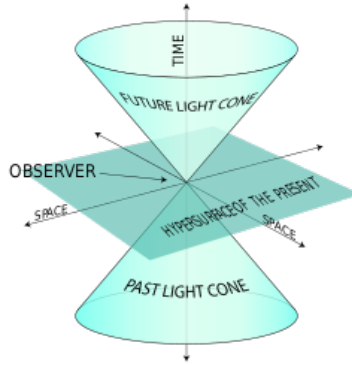


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So we need to be more specific on *where* we measure things, since measurements separated by space-like vectors should not interfere with each other.

We need to take into account special relativity: the Minkowski geometry of space–time and the limitation on the speed of allowed interactions among different regions of space–time. Accordingly our measurement algebra need to encode the fact that space-like separated experiments should not interfere with each other, i.e. the associated algebras commute (one should say this more precisely but we will skip over these fine points). Naively one could try to associate an observable algebra to each point of a equal time surface (because they are all space-like with respect to each other, i.e. outside the light-cone of each other). A more conservative approach is to consider each measuring apparatus to have a finite spatial precision and measure what happens in an extended (even if maybe small) region of space. More abstractly we could associate to any compactly supported smooth test function $f \in \mathcal{D}(\mathbb{R}^n)$ in n -dimensional space a one-dimensional observable algebra (for example), i.e. an homeomorphism $\Phi_0(f): C(\mathbb{R}; \mathbb{C}) \rightarrow \mathcal{A}$ into the C^* -algebra \mathcal{A} generated by all observables. By writing the more suggestive notation $F(\Phi_0(f)) := \Phi_0(f)(F)$ for any $F \in C(\mathbb{R}; \mathbb{C})$ we would like to impose that

$$[F(\Phi_0(f)), G(\Phi_0(g))] = 0$$

for each $f, g \in \mathcal{D}(\mathbb{R}^n)$ with disjoint support and $F, G \in C(\mathbb{R}; \mathbb{C})$.

All together we then expect a representation of the full Poincaré group by automorphisms of \mathcal{A} .

For example if we assume that we can make measurements at a precise point, we can assume to have observables $F(\Phi(t, x))$ with measures a given quantity at time $t \in \mathbb{R}$ and in the point $x \in \mathbb{R}^n$, $(t, x) \in \mathbb{M}^n$. The an element g of the Poincaré group (i.e. the symmetry group of \mathbb{M}^n) acts via an automorphism α_g such that

$$\alpha_g(F(\Phi(t, x))) = F(\Phi(g.(t, x))), \quad (t, x) \in \mathbb{M}^{n+1}$$

where $g.(t, x) \in \mathbb{M}^n$ denotes the action of g on the vector (t, x) . Moreover $\alpha_g \circ \alpha_{g'} = \alpha_{gg'}$. (Implicitly I'm considering a scalar field).

Accordingly is also natural to assume that there exists a state ω which is invariant under all the $(\alpha_g)_g$ and that the group $(\alpha_g)_g$ has some mild regularity properties.

With some additional mild assumptions one can then derive the existence of a Hilbert space \mathcal{H} , a state vector φ , an strongly continuous unitary representation $(U_g)_g$ of the Poincaré group which leaves the state vector φ invariant and a representation $\pi \in \text{Hom}(\mathcal{A}, \mathcal{B}(\mathcal{H}))$ of the observable algebra on \mathcal{H} .

Note that $[F(\Phi(t, x)), G(\Phi(t, y))] = 0$ for all F, G and all $x \neq y$. So I can consider all the collection of all $(\Phi(0, x))_{x \in \mathbb{R}^n}$ for all the points at time zero. They form a commutative subalgebra of observables.

Let us now apply the Euclidean strategy and look for a measure which could deliver such data.

We expect the stochastic process X to be, at each point of time $t \in \mathbb{R}$ a random distribution, i.e. that for each $f \in \mathcal{D}(\mathbb{R}^n)$ we can construct a random variable $X_t(f)$ and

$$X_t(\alpha f + \beta g) = \alpha X_t(f) + \beta X_t(g).$$

That is: we have an infinite-dimensional random process, apart from this nothing much changes to the previous finite dimensional considerations above, in particular RP works just fine.

(However as we mentioned above the Markovian strategy is more subtle in infinite dimensions (i.e. conditioning wrt. \mathcal{F}_0 is a much stronger perturbation of the system and it is not clear that it is equivalent to just conditioning on X_0) and we adopt the strategy of perturbing a simple RP process.)

The goal is the to find RP processes in space-time which are invariant wrt. to time and space translations. (What about Poincaré rotations, i.e. rotations and boosts?) One realises that the Poincaré invariance corresponds to the Euclidean invariance of the stochastic process in $\mathbb{R}^{n+1} \approx \mathbb{R}^d$. That is we want the law of X to be invariant under the full Euclidean group acting on $(t, x) \in \mathbb{R}^d$.

► This justifies our definition of EQFT: i.e. probability measure on distributions over \mathbb{R}^d that is RP and Euclidean invariant (+ other technical regularity assumptions).

Let us look therefore at what is available. We can try to take a process X such that the family $(X_t(f))_{t \in \mathbb{R}, f \in \mathcal{D}(\mathbb{R}^n)}$ is jointly Gaussian and we know that to get a RP process we can look at OU processes.

Let us start assuming that we build $X_t(f)$ as a sum of finitely many independent OU (complex) processes (Y_t^i) as follows

$$X_t(f) = \sum_i \hat{f}(k_i) Y_t^i, \quad (6)$$

where \hat{f} is the Fourier transform of f and $k_i \in \mathbb{R}^n$ are parameters. Precisely we take

$$X_t(f) = \sum_i \operatorname{Re}(\hat{f}(k_i)) Y_t^{1,i} + \operatorname{Im}(\hat{f}(k_i)) Y_t^{2,i}.$$

Then we have to arrange the covariances such that

$$\begin{aligned} \mathbb{E}[X_t(f) X_0(f)] &= \sum_{i,j} \operatorname{Re}(\hat{f}(k_i)) \operatorname{Re}(\hat{f}(k_j)) \underbrace{\mathbb{E}[Y_t^{1,i} Y_t^{1,j}]}_{=\delta_{i,j} c_i e^{-\lambda_i |t|}} + \sum_{i,j} \operatorname{Im}(\hat{f}(k_i)) \operatorname{Im}(\hat{f}(k_j)) \underbrace{\mathbb{E}[Y_t^{2,i} Y_t^{2,j}]}_{=\delta_{i,j} c_i e^{-\lambda_i |t|}} \\ &+ \sum_{i,j} \operatorname{Re}(\hat{f}(k_i)) \operatorname{Im}(\hat{f}(k_j)) \underbrace{\mathbb{E}[Y_t^{1,i} Y_t^{2,j}]}_{=0} + \sum_{i,j} \operatorname{Im}(\hat{f}(k_i)) \operatorname{Re}(\hat{f}(k_j)) \underbrace{\mathbb{E}[Y_t^{2,i} Y_t^{1,j}]}_{=0} \end{aligned}$$

Then the covariance reads:

$$\mathbb{E}[X_t(f) X_0(f)] = \sum_i |\hat{f}(k_i)|^2 c_i e^{-\lambda_i |t|}.$$

Exercise. make this precise.

Note that

$$e^{-\lambda_i |t|} = \pi \int_{\mathbb{R}} \frac{e^{i\omega t}}{\omega^2 + \lambda_i^2} d\omega.$$

(modulo the right coefficient). So

$$\mathbb{E}[X_t(f)X_0(f)] \propto \int_{\mathbb{R}} \sum_i |\hat{f}(k_i)|^2 c_i \frac{e^{i\omega t}}{\omega^2 + \lambda_i^2} d\omega$$

Now taking more and more points and appropriate c_i we can enforce convergence to a integral over $k \in \mathbb{R}^n$:

$$\mathbb{E}[X_t(f)X_0(f)] = \int_{\mathbb{R}^n} \int_{\mathbb{R}} |\hat{f}(k)|^2 \frac{e^{i\omega t}}{\omega^2 + \lambda(k)^2} dk d\omega$$

(we would need now infinitely many degrees of freedom, i.e. OU processes).

This expression implies that the spatial covariance has the form

$$\mathbb{E}[X_t(f(\cdot+x))X_0(f)] = \int_{\mathbb{R}^n} \int_{\mathbb{R}} |\hat{f}(k)|^2 \frac{e^{i(\omega t+k \cdot x)}}{\omega^2 + \lambda(k)^2} dk d\omega$$

so it makes sense to take $\lambda(k)^2 = k^2 + m^2$ for a fixed constant $m > 0$. This give rise to an Euclidean invariant covariance. To see this take $f \rightarrow \delta$ so $\hat{f}(k) \rightarrow 1$ and observe that one obtains

$$\mathbb{E}[X_t(\delta_x)X_0(\delta_0)] = \int_{\mathbb{R}^d} \frac{e^{ik \cdot (t,x)}}{k^2 + m^2} dk.$$

The associated process is therefore Euclidean invariant in \mathbb{R}^d , since the covariance is. This form of the covariance is constrained by RP and Euclidean invariance.

Definition. *This process is called the Gaussian Free field with mass $m > 0$ in (Euclidean) dimension d . This this an example of EQFT.*

For $d = 1$ is just the OU process we have seen.

Note that we have problems... :

$$\mathbb{E}[|X_0(\delta_0)|^2] = \int_{\mathbb{R}^d} \frac{1}{k^2 + m^2} dk = +\infty,$$

if $d \geq 2$. Note that $d = 1$ is just QM... (finitely many degrees of freedom). I cannot evaluate X at any fixed space-time point. It is still true that for any $f \in \mathcal{D}(\mathbb{R}^n)$ the process $(X_t(f))_{t \in \mathbb{R}}$ is a nice and continuous Gaussian process.

From eq. (6) we have

$$dX_t(f) = \sum_i \hat{f}(k_i) dY_t^i = \sum_i \hat{f}(k_i) \lambda_i Y_t^i dt + \sum_i \hat{f}(k_i) dB_t^i = X_t((m^2 - \Delta)^{1/2} f) dt + dB_t(f)$$

with

$$B_t(f) = \sum_i \hat{f}(k_i) B_t^i.$$

In the limit we obtain therefore the Euclidean-time dynamics

$$dX_t(f) = X_t((m^2 - \Delta)^{1/2} f) dt + dB_t(f)$$

with $B_t(f)$ a Brownian motion such that $\mathbb{E}[B_t(f)^2] = t \|f\|_{L^2(\mathbb{R}^n)}^2$.

Note that our construction guarantees that it is reflection positive in the time direction. By rotation invariant we have also that it is reflection positive along any direction.

Full Euclidean invariance is important because, within the reconstruction of QM guarantees that the quantum theory is Poincaré invariant.

Exercise. Try to work out some details of the reconstruction in this case, it would be nice to arrive at a point where one can check explicitly the Poincaré invariance of the quantum theory.

3.5 Interacting Euclidean Quantum Fields?

From now on we stop to treat the physical time as a special variable and consider it part of the Euclidean coordinates and we will work in \mathbb{R}^d forgetting about the Minkowski geometry. This is certainly one advantage of the Euclidean approach. As we have shown in the one dimensional setting we can go back to a full quantum theory. Here again the theory is trivial: the dynamics do not mix the different OU processes. This also implies that one can factorize the quantum Hilbert space into a product of independent components, each of them behaving as we described in one dimension: the system is full of non-interacting particles.

To construct other examples out of the GFF one can use the perturbative approach and consider an interaction potential of the form $V(X_t)$. However doing so we break the Euclidean invariance of the theory: one one hand we need that the Gibbsian perturbation has the form

$$\exp\left(\int_{-T}^T V(X_s) ds\right)$$

to maintain RP, on the other hand this special treatment of the time direction should be compensated by a specific form of V to reestablish some invariance under rotations. We are led to take the perturbation to be in the form

$$\exp\left(\int_{[-T,T]^d} V(X(x)) dx\right)$$

where we denote $X(x_1, x_2, \dots, x_d) = X_{x_1}(x_2, \dots, x_d)$ the OU process and consider the new measure

$$d\mathbb{Q}^T = \frac{1}{Z_T} \exp\left(\int_{[-T,T]^d} V(X(x)) dx\right) d\mathbb{P} \quad (7)$$

with suitable normalization constant.

Now we encounter more difficulties: first of all we had to cut of the integral in all the dimensions to have a well defined quantity: we call T an infrared cutoff: it limits the perturbation to a compact region of space, giving hope to gave a meaning to the perturbed measure and thus breaking even more the translation invariance of the model. As with the finite dimensional situation however we can hope that the $T \rightarrow \infty$ limit reestablish such invariance. This is also confirmed by statistical mechanical models which have very strong similarities with the measure (?) and for which one understand well what happens when $T \rightarrow \infty$.

A more serious problem is due to UV (ultraviolet) divergences. As soon as one stops a moment and think, one realises that $X(x)$ is not a well defined random variable, indeed we have

$$\mathbb{E}[X(f)^2] = \int_{\mathbb{R}^d} \frac{|f(k)|^2}{k^2 + m^2} dk.$$

so as $f \rightarrow \delta_x$ this variance explodes if $d \geq 2$ due to failure of the integrability at infinity caused by the slow decay of the integrand (whose specific form was motivated by Euclidean invariance). Recall that the case $d = 1$ corresponds to the one dimensional OU process, so of no interest of us at this point.

We are forced to admit the possibility that the samples of the random variable X are not regular enough to be evaluated at points $x \in \mathbb{R}^d$ and we have to resort to a more careful handling. In particular the expression $V(X(x))$ does not make sense. This is a problem due to the fluctuations at small scales so it is called an UV divergence (because UV light rays oscillates at smaller scales than IR light rays).

To have a sensible starting point we therefore add another layer of approximation and avoid very large fluctuations. This can be achieved in many ways. An easy approach is to fix a smooth function $\rho: \mathbb{R}^d \rightarrow \mathbb{R}_+$ so that $\int \rho(x) dx = 1$ and with compact support, let $\rho_\varepsilon(x) = \varepsilon^{-d} \rho(x/\varepsilon)$ to rescale it while preserving the L^1 norm and then let $X_\varepsilon = \rho_\varepsilon * X_\varepsilon$ which can be checked to be a nice and smooth gaussian process, that is we can prove that \mathbb{P} -almost surely $X_\varepsilon \in C^\infty(\mathbb{R}^d)$ (it is not bounded however and the control of the smoothness degrades with the distance from the origin). We will come later on to details on these considerations, for the moment let's assume this fact. Keeping $\varepsilon > 0$ we can now consider

$$d\mathbb{Q}^{\varepsilon, T} = \frac{1}{Z_T} \exp\left(\int_{[-T, T]^d} V(X_\varepsilon(x)) dx\right) d\mathbb{P} \quad (8)$$

which is well defined, e.g. when V is bounded below and satisfying some additional mild assumptions. For example we can take $V(\varphi) = \varphi^4$ or any other polynomial bounded from below and things would be fine.

We have arrived finally to a meaningful and nontrivial mathematical problem: prove that the family of measures $\mathbb{Q}^{\varepsilon, T}$ converges weakly to a measure \mathbb{Q} on random distributions on \mathbb{R}^d .

Of course this is not enough to carry on the reconstruction program: one needs that the limit satisfies RP, stationarity and an additional analytic condition.

If moreover the measure \mathbb{Q} is also rotation invariant, then the reconstructed QM is Poincaré invariant.

The above choice of UV regularisation destroys completely RP and it is not clear how to recover this property in the limit. Another possible regularisation is to apply the smoothening with ρ_ε in all directions but one, for which then the proof of RP given above works.

In these notes we will follow a technically easier approach to construct first a model on a lattice instead of continuous space, i.e. we replace $[-T, T]^d$ by $\mathbb{T}_{L, \varepsilon}^d := [-L, L]^d \cap (\varepsilon\mathbb{Z})^d$ a discrete set of $(L/\varepsilon)^d$ points which we endow with periodic boundary conditions. The boundary conditions are useful to have the translation invariance of the model, since the GFF is not periodic of period L we need also to periodicize it. The result is that we will consider a family of Gaussian r.v. $(X^{L, \varepsilon}(x))_{x \in \mathbb{T}_{L, \varepsilon}^d}$ with covariance

$$\mathbb{E}[X^{L, \varepsilon}(x) X^{L, \varepsilon}(y)] = (m^2 - \Delta_\varepsilon)^{-1}(x, y),$$

where Δ_ε is the discrete Laplacian with periodic boundary conditions. A Fourier transform formula for this function reads

$$\mathbb{E}[X^{L, \varepsilon}(x) X^{L, \varepsilon}(y)] = \frac{1}{(2\pi L)^d} \sum_{k \in ((\mathbb{Z}/L) \cap [-\varepsilon^{-1}, \varepsilon])^d} \frac{e^{ik \cdot (x-y)}}{(\sin(\varepsilon k) \varepsilon^{-1})^2 + m^2}, \quad x, y \in [-L, L]^d \cap (\varepsilon\mathbb{Z})^d.$$

The exact formula is not very important.

Then we take $\mu^{L, \varepsilon}$ to be the law of $X^{L, \varepsilon}$ and define $\nu^{L, \varepsilon}$ as

$$d\nu^{L, \varepsilon} = \frac{1}{Z_{L, \varepsilon}} \exp\left(-\varepsilon^d \sum_{x \in \mathbb{T}_{L, \varepsilon}^d} V(X^{L, \varepsilon}(x))\right) d\mu^{L, \varepsilon} \quad (9)$$

which is now a very well defined (even elementary) expression.

The advantage of this discretisation is that it is invariant under lattice translations and that satisfy a discrete version of reflection positivity [see. e.g. Sacha Friedli and Yvan Velenik, *Statistical Mechanics of Lattice Systems: A Concrete Mathematical Introduction* (Cambridge, United Kingdom; New York, NY: Cambridge University Press, 2017).].

This will be our starting point to discuss stochastic quantisation.