

# **Perturbative Quantum Field Theory at Positive Temperatures: An Axiomatic Approach**

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**Abstract:** It is shown that the perturbative expansions of the correlation functions of a relativistic quantum field theory at finite temperature are uniquely determined by the equations of motion and standard axiomatic requirements, including the KMS condition. An explicit expression as a sum over generalized Feynman graphs is derived. The canonical formalism is not used, and the derivation proceeds from the beginning in the thermodynamic limit. No doubling of fields is invoked. An unsolved problem concerning existence of these perturbative expressions is pointed out.

## **1. Introduction**

The traditional way of describing thermal equilibrium states of an infinitely extended quantum system, in particular of a quantum field theory, begins by restricting the system to a finite volume  $V$ , defining the canonical or grand canonical equilibrium by means of the familiar density matrices, and then going to the limit  $V \rightarrow \infty$  (the “thermodynamic limit”) for the quantities for which this limit can be expected to exist [1, 2]. This applies especially to the correlation functions of the fields and closely related objects like the expectation values of time ordered field products. Up to now most actual calculations of such functions have been based on this approach, using a Hamiltonian or Lagrangian formalism at finite  $V$ .

Another description of equilibria and their local disturbances, which can be used directly in the thermodynamic limit, has been developed in the framework of the algebras of local observables [3]. In this approach equilibrium states are characterized through an analyticity requirement for correlation functions, the so-called KMS condition. In the present paper we intend to show that this axiomatic method, suitably adapted to a field theoretical context, is perfectly capable of handling dynamical problems. More exactly, it will be shown that perturbative expansions for the correlation functions of a relativistic field theory, and related functions, can be derived directly in the thermodynamic limit, not making use of the canonical formalism, but using as only inputs the equations of motion and the axiomatic requirements that the correlation functions must satisfy. The result is represented as a sum over

generalized Feynman graphs. For the special case of time ordered functions it agrees with the well-known result of the canonical approach.

Dispensing with the canonical formalism is also a major difference between our approach and thermo field dynamics [2, 4], a Fock space method developed by H. Umezawa and coworkers. We differ also from thermo field dynamics by not invoking a doubling of fields, and by not assigning a basic role to particles, including quasi-particles. We hold particles to be secondary objects of the theory, of great phenomenological importance, but little fundamental significance. In this respect we differ also from the views put forward by Landsman in ref. [5].

We consider only the  $\Phi_4^4$ -model, i.e. the theory of a scalar hermitian field  $\Phi(x)$  satisfying the equation of motion

$$(\square + m^2)\Phi = -\frac{g}{6}N(\Phi^3), \quad (1.1)$$

where  $N$  stands for a normal-product prescription taking care of renormalization. The restriction to  $\Phi_4^4$  is merely a matter of convenience. The generalization of the method to other models, including gauge theories in local gauges, is straightforward.

We are interested in the correlation functions

$$W(x_1, \dots, x_n) = \langle \Phi(x_1) \dots \Phi(x_n) \rangle, \quad (1.2)$$

where  $\langle \cdot \rangle$  denotes the expectation value in a thermal equilibrium state with temperature  $T \geq 0$ . These correlation functions describe the full physical content of the theory: all observable quantities can in principle be derived from them (for examples see e.g. [1, 2]). This is so because knowledge of the  $W$  allows the reconstruction of the full representation of the field algebra by means of the GNS construction [2, 3], which yields a Hilbert space representation of the field algebra with a cyclic vector  $|\cdot\rangle$ , such that  $\langle A \rangle = \langle |R(A)| \rangle$  for any sufficiently regular function  $A(\Phi)$ , with  $R(A)$  its representative. All observables of the theory are supposed to be of this form, and local disturbances of the equilibrium are created by applying suitable functions  $F(\Phi)$  to  $|\cdot\rangle$ .

More generally, we consider the set of functions (or rather, distributions)

$$\mathcal{W}(X_1, s_1 | \dots | X_N, s_N) = \langle T^{s_1}(X_1) \dots T^{s_N}(X_N) \rangle. \quad (1.3)$$

Here the  $X_x$  are non-overlapping sets of 4-vectors  $x_i$ , the  $s_x$  are signs, and  $T^\pm(X)$  denotes respectively the time-ordered or anti-time-ordered product of the fields  $\Phi(x_i)$ ,  $x_i \in X$ . If each  $X_x$  contains only one variable, then  $\mathcal{W}$  is the correlation function  $W(x_1, \dots, x_N)$  irrespective of the choice of the signs  $s_x$ . For  $N = 1$  we obtain the usual time-ordered and anti-time-ordered functions (Green's functions) of the theory. In the sequel the signs  $s_x$  will be frequently suppressed when they are not essential to understanding.

In a previous work [6], henceforth quoted as V, we have derived perturbative expressions for the functions  $\mathcal{W}$  in terms of generalized Feynman graphs for the case  $T = 0$ , in which case  $\langle \cdot \rangle$  denotes the vacuum expectation value. This derivation uses neither the Hamiltonian nor the Lagrangian formalism, but is instead relying on the Wightman axioms as an essential input. We propose to generalize this method to the case  $T > 0$ . Our way of proceeding is closely modeled on that taken in V. The ideas and results of that paper will be freely used. Equation (n.m.) of V will be referred to as Eq. (V.n.m.). The method consists in solving the differential equations

for  $W$  which follow from the field equation (1.1), by a power series expansion in the coupling constant  $g$ , using the axiomatic properties of  $W$  as subsidiary conditions.

The paper is organized as follows. The assumptions on which the formalism is based will be stated and discussed in Sect. 2. An explicit formal expression for the  $\mathcal{W}$ 's as sums over generalized Feynman graphs will be stated in Sect. 3 and shown to possess the required properties. In Sect. 4 the ultraviolet (UV) divergences of these graphs will be removed by renormalization. It will be pointed out, however, that renormalization does not guarantee the existence of the resulting finite-order expressions, on account of the local singularities of the integrands. This problem remains unsolved. Finally we will show in Sect. 5 that the expressions of Sects. 3 and 4, provided they exist, are the only ones satisfying the assumptions stated in Sect. 2.

## 2. Assumptions

In this section the assumptions on which our derivations are based, in addition to the field equation (1.1), will be enumerated.

The following conditions for  $W$  and  $\mathcal{W}$  are taken over unchanged from V:

a) The  $W$  and  $\mathcal{W}$  are *invariant* under space-time translations and under space rotations. Invariance under Lorentz boosts cannot be demanded for  $T > 0$ .

b) *Locality* holds, i.e.  $W(X)$  is invariant under the exchange of two neighbouring variables  $x_i, x_{i+1}$ , if  $x_i - x_{i+1}$  is space-like.

c) The *reality condition*

$$\mathcal{W}(X_1, s_1 | \dots | X_N, s_N)^* = \mathcal{W}(X_N, -s_N | \dots | X_1, -s_1) \tag{2.1}$$

holds.

d) The functions  $\mathcal{W}$  are *permutation invariant* within each sector  $X_\alpha$ , and they satisfy the *splitting property*

$$\mathcal{W}(\dots | X_1 \cup X_2, + | \dots) = \mathcal{W}(\dots | X_1, + | X_2, + | \dots) \tag{2.2}$$

if  $x_i^0 > x_j^0$  for all  $x_i \in X_1, x_j \in X_2$ . This property shall hold in every Lorentz frame, not only in the rest frame of the infinite system under consideration. To that extent we retain Lorentz invariance. These conditions are not merely a definition of time ordering. That they can be satisfied is intimately connected with locality.

The *cluster property* needs a more careful discussion than it was accorded in V. Let  $X, Y$ , be two non-overlapping sets of 4-vectors. Then the cluster property states that

$$\lim_{|a| \rightarrow \infty} W(X, Y + a) = W(X)W(Y), \tag{2.3}$$

where  $Y + a$  means that all vectors in  $Y$  are translated by  $a$ , and  $a$  tends to infinity in a space-like direction. Since derivation with respect to the coupling constant need not commute with the limit in (2.3), we cannot expect Eq. (2.3) to hold separately in each order of perturbation theory, except in the lowest, free, order. But perturbation theory can also be viewed as an expansion in powers of  $\hbar$ . And we can demand that Eq. (2.3) hold for each  $W(X)$  in the lowest nonvanishing order in  $\hbar$ . This is all that will be needed to establish uniqueness. We set  $c = 1$ , so that time and space have the same dimension. The dimension of the field  $\Phi$  is then  $[\text{mass}/\text{length}]^{1/2}$ . In order

to make the field equation (1.1) dimensionally consistent, the coupling constant in the interaction term reads  $g/\hbar$ , with a dimensionless  $g$ . We shall nevertheless set  $\hbar = 1$  in our equations, and merely note the correct  $\hbar$ -exponents at the points where they are essential.

The spectral condition of the vacuum representation does not hold at positive temperatures. It is replaced by the *KMS condition*, which we use in its p-space form (see [3], Lemma 1.1.1 of Chap. V). Define the Fourier transform  $\tilde{\Phi}(p)$  of  $\Phi$  by

$$\tilde{\Phi}(p) = (2\pi)^{-5/2} \int dx e^{ipx} \Phi(x), \tag{2.4}$$

and the Fourier transform  $\tilde{\mathcal{W}}(P_1 | \dots | P_N)$  of  $\mathcal{W}(X_1 | \dots | X_N)$  accordingly. Let  $P_1, \dots, P_N, Q_{N+1}, \dots, Q_M, M > N$ , be finite, non-empty sets of 4-momenta. Define

$$P^0 = \sum_{p_i \in \cup P_x} p_i^0. \tag{2.5}$$

Then the KMS condition states that

$$\tilde{\mathcal{W}}(P_1 | \dots | P_N | Q_{N+1} | \dots | Q_M) = e^{\beta P^0} \tilde{\mathcal{W}}(Q_{N+1} | \dots | Q_M | P_1 | \dots | P_N), \tag{2.6}$$

where  $\beta = \frac{1}{kT}$  is the inverse temperature.

The normalization conditions stipulated in V are taken over unchanged for the case  $T = 0$ . They are the standard conditions demanding that  $m$  be the physical mass of the particles associated with the field  $\Phi$  at  $T = 0$ , and defining the coupling constant  $g$  and the field normalization in terms of the Green's functions. In addition we demand that  $\Phi$ , considered as an element of an abstract, representation independent, field algebra, is  $T$ -independent. This implies that the field equation is  $T$ -independent, meaning that the parameters  $m$  and  $g$  and the subtraction prescription  $N$  do not depend on the temperature. The prescription  $N$  is assumed to be the conventional one. In particular it should not contain oversubtractions which would destroy the renormalizability of the theory. In perturbation theory the  $T$ -independence of  $N$  means that the UV subtractions and the subsequent finite renormalizations have the values used at  $T = 0$  for all  $T$ . Note that in the BPHZ method [7], which we will be using, the subtractions do not involve prior regularizations or any formal juggling with divergent quantities. The procedure is therefore well defined. For more details we refer to Sect. 4. A nonperturbative definition of  $N$  can possibly be given with the help of a Wilson expansion, as explained in Sect. IV.2. of ref. [7], i.e. by a point-splitting method with suitable singular coefficient functions.  $T$ -independence of  $N$  could then be defined as  $T$ -independence of these coefficients. We will not explore this possibility any further.

Note that the parameter  $m$  is for  $T > 0$  not the physical mass of a particle. Quasi-particle masses are determined by the positions of the singularities of the clothed propagator, which are  $T$ -dependent.

### 3. Unrenormalized Solution

The coefficient  $W_\sigma(X)$  of order  $\sigma$  in the power series expansion

$$W(X) = \sum_{\sigma=0}^{\infty} g^\sigma W_\sigma(X) \tag{3.1}$$

is determined as a solution of the system of differential equations

$$(\square_i + m^2)\langle \dots \Phi(x_i) \dots \rangle_\sigma = -\frac{1}{6}\langle \dots N(\Phi(x_i)^3) \dots \rangle_{\sigma-1} \tag{3.2}$$

satisfying the subsidiary conditions described in Sect. 2. For  $\sigma = 0$  the right-hand side is zero, for  $\sigma > 0$  it can be calculated if the problem has already been solved in order  $\sigma - 1$ .

We first state a formal, unrenormalized, graphical representation for  $\mathcal{W}_\sigma(X_1 | \dots | X_N)$  and then show that it does indeed satisfy all the requirements.

For the free propagators we use the following notations. We define first the p-space expressions

$$\tilde{D}_+(p) = \frac{1}{(2\pi)^3} \delta_+(p), \quad \tilde{D}_F(p) = \frac{i}{(2\pi)^4} \frac{1}{p^2 - m^2 + i\varepsilon}, \tag{3.3}$$

which are the usual vacuum propagators, and their thermal extensions

$$\tilde{D}_+(p) = \tilde{D}_+(p) + \tilde{C}(p), \quad \tilde{D}_F(p) = \tilde{D}_F(p) + \tilde{C}(p) \tag{3.4}$$

with the  $T$ -dependent correction term

$$\tilde{C}(p) = \frac{1}{(2\pi)^3} \frac{1}{e^{\beta\omega} - 1} \delta(p^2 - m^2). \tag{3.5}$$

Here  $\delta_+(p) = \theta(p_0)\delta(p^2 - m^2)$  is the  $\delta$ -function of the positive mass shell, and  $\omega = (\mathbf{p}^2 + m^2)^{1/2}$ . Note that the additional term  $\tilde{C}$  is the same for  $\tilde{D}_+$  and  $\tilde{D}_F$ .

The  $x$ -space versions of these propagators are defined by

$$\Delta \cdot (x) = i \int d^4 p \tilde{A} \cdot (p) e^{-ipx}, \quad D \cdot (x) = i \int d^4 p \tilde{D} \cdot (p) e^{-ipx}. \tag{3.6}$$

The thermal propagators (3.4-5) need not be taken over from the traditional formalism. They can be derived in our framework by the methods that will be used in Sect. 5 to establish uniqueness of our solution. But in the present section we rely on an a posteriori justification, by showing that these forms give rise to expressions with all the required properties.

$\mathcal{W}_\sigma$  is represented as a sum over generalized Feynman graphs which are defined as follows. Draw first an ordinary Feynman graph of the  $\Phi^4$  theory with  $|X| = \sum_x |X_x|$  external and  $\sigma$  internal vertices. Here  $|X_x|$  is the number of points in the set  $X_x$ . The graph need not be connected, but must not contain any components without external points. This graph is called the “scaffolding” of the generalized graph. Next, it is partitioned into non-overlapping subgraphs, called “sectors,” such that the external points of a set  $X_x$  belong all to the same sector, but variables in different  $X_x$  to different sectors. There may also exist “internal sectors” not containing external points. The sectors are either of type  $T^+$  or  $T^-$ . For external sectors this sign is given by  $s_x$ . To each sector  $S$  we affix its number  $\nu(S)$  according to the following rules:

- i)  $\nu(S) = \alpha$  for the external  $X_x$ -sector.
- ii) If  $s_x = s_{x+1}$  there may be an internal sector with number  $\alpha + \frac{1}{2}$ . Its type is the reverse of the adjacent external sectors:  $s_{x+\frac{1}{2}} = -s_x$ . If  $s_x = -s_{x+1}$  no such intermediate internal sector exists.

- iii) If  $s_1 = s_N$  there may be an internal sector with number  $N + \frac{1}{2}$  and  $s_{N+\frac{1}{2}} = -s_N$ . Equivalently we could give this additional sector the number  $\frac{1}{2}$ , but we must choose one of these possibilities and use it consistently.

With a partitioned graph we associate a Feynman integrand, which we state first in  $x$ -space. To the external points correspond the external variables  $x_i$ . To each internal vertex we assign an integration variable  $u_j$ ,  $j = 1, \dots, \sigma$ .  $z_i$  denotes a variable which may be either external or internal. Within a  $T^+$  sector the usual Feynman rules hold: each internal vertex carries a factor  $-ig$ , a line connecting the points  $z_i$  and  $z_j$  carries the propagator  $-iD_F(z_i - z_j)$ . In a  $T^-$  sector the complex-conjugates of these rules apply. A line connecting points  $z_i, z_j$  in different sectors, with  $z_i$  lying in the lower-numbered sector, carries the propagator  $-iD_+(z_i - z_j)$ . The graph must be divided by the usual symmetry number if it is invariant under certain permutations of points and lines.

In  $p$ -space, integration variables are assigned to the lines. The vertex factors in  $T^\pm$  factors are  $\mp i(2\pi)^4 g$  and the propagators are  $\tilde{D}_F(p)$  or  $(\tilde{D}_F(p))^*$  respectively. Lines connecting different sectors carry propagators  $\tilde{D}_+(p)$ , the momentum  $p$  flowing from the lower to the higher sector. For each external point there is a factor  $(2\pi)^{3/2}$ . Momentum is conserved at each internal vertex.

For  $T = 0$  the rule ii) governing internal sectors seems to differ from the corresponding rule in  $V$ , where instead of one (possibly empty) intermediate sector of different type we had chains of intermediate sectors of the same type as the bracketing external sectors. However, it can be shown with the help of Lemma V.3.1 that the two formulations are equivalent. Graphs containing non-empty internal sectors according to iii) vanish for  $T = 0$ , wherefore they do not occur in  $V$ .

For the Green's functions ( $N = 1$ ) our rules agree with those of the conventional real-time formalism. Our graphs are equal to those of the Keldysh formulation [8].

The first point to be checked is that the above prescription gives an unambiguous result for the correlation functions  $W$ . The problem is that in the definition  $W(x_1, \dots, x_n) = \mathcal{W}(x_1, s_1 | \dots | x_n, s_n)$  the signs  $s_x$  can be chosen arbitrarily. But it can be shown that the sum over all graphs with the same scaffolding does not depend on the choice of these signs. The proof is modeled closely on the corresponding proof in  $V$  and will only be briefly indicated. Consider a given scaffolding. Let  $S$  denote an internal  $T^+$  sector,  $S_x$  an external  $T^+$  sector with external variable  $x, \bar{S}$  and  $\bar{S}_x$   $T^-$  sectors. A product  $S_x S$  or the like denotes the sum over all partitions of a given subgraph into two adjacent sectors  $S_x, S$ . The propagators connecting these two sectors are included. We will first show that the choice of  $s_N$  is irrelevant. The graph's belonging to  $s_N = +$  or  $s_N = -$  respectively differ only in the sectors with number  $v(S) > N - 1$ . In the case  $s_{N-1} = +, s_1 = -$  these variable sectors are (we put  $x_N = x$ )  $S_x + \bar{S}S_x$  for  $s_N = +$ , and  $\bar{S}_x + \bar{S}_x S$  for  $s_N = -$ . But

$$S_x + \bar{S}S_x = \bar{S}_x + \bar{S}_x S$$

is a special case of Lemma V.3.1, which remains valid for our new propagators. Similarly, if  $s_1 = s_{N-1} = +$  we must show that

$$S_x + S_x \bar{S} + \bar{S}S_x + \bar{S}S_x \bar{S} = \bar{S}_x,$$

or

$$(S_x + S_x \bar{S} - \bar{S}_x - \bar{S}S_x) + \bar{S}(S_x + S_x \bar{S} - \bar{S}_x - \bar{S}S_x) + (S + \bar{S} + \bar{S}S)\bar{S}_x = 0,$$

which is correct because all three brackets vanish as a result of Lemma V.3.1. The proof of independence on the other  $s_x$  follows similar lines. The case  $1 < \alpha < N$  is considerably simplified compared to V by the new formulation of the rules concerning internal sectors.

Of the conditions stated in Sect. 2 invariance under translations and rotations, and the symmetry of  $\mathcal{W}$  within  $T^\pm$  factors are trivially satisfied. That the equation of motion (3.2) is satisfied, is shown exactly as in V. At the present formal level the product  $N(\Phi^3)$  denotes Wick ordering  $:\Phi^3(x):$  with respect to the free vacuum field.

The proof of the reality condition (2.1) and Ostendorf's proof [9] of the splitting relation (2.2) can also be taken over unchanged from V. Since the non-invariant part  $C(x - y)$  is the same in  $D_+$  and in  $D_F$ , the crucial relation  $D_F(x - y) = D_+(x - y)$  if  $x^0 > y^0$  holds in every orthochronous Lorentz frame. Hence the same is true for the splitting relation. From this and the symmetry within sectors we can prove locality. Let  $(x - y)^2 < 0$ . Then there exist two Lorentz frames such that in one  $x^0 > y^0$ , in the other  $x^0 < y^0$ . Hence

$$\begin{aligned} \mathcal{W}_\sigma(\dots, x, y, \dots) &= \mathcal{W}_\sigma(\dots |x, y, + | \dots) \\ &= \mathcal{W}_\sigma(\dots |y, x, + | \dots) = \mathcal{W}_\sigma(\dots, y, x, \dots). \end{aligned}$$

Our propagators, being the 2-point functions of a free field, contain for dimensional reasons a factor  $\hbar$  (this refers to x-space). The vertices carry a factor  $\hbar^{-1}$ , so that the  $2n$ -point function of order  $\sigma$  contains a factor  $\hbar^{n+\sigma}$ . The term of lowest order in  $\hbar$  is therefore the free  $\sigma = 0$  expression. It is simply a sum over all possible products of free 2-point functions, and obviously satisfies the cluster property, because  $D_+(\xi)$  and  $D_F(\xi)$  converge to zero if  $\xi$  tends to infinity in a space-like direction. Hence our weak version of the cluster property is satisfied.

Yet to be proved remains the KMS condition (2.6). It is easy to see that a graph contributing to the left-hand side of (2.6) becomes a graph contributing to the right-hand side, if the directions of all lines connecting sectors with  $\nu(S) < N + 1$  to sectors with  $\nu(S) \geq N + 1$  are reversed, and vice versa. This means that the corresponding propagators  $\tilde{D}_+(k)$  are replaced by  $\tilde{D}_+(-k)$ . From the definitions (3.3–5) we find  $\tilde{D}_+(k) = e^{\beta k_0} \tilde{D}_+(-k)$ . Hence the two variants of the considered graph differ by the factor  $\exp(\beta \sum_x k_0^x)$ , the sum extending over the momenta of the lines in question. But  $K = \sum k^x$  is the total momentum flowing from the  $P$ -part of the graph to the  $Q$ -part, so that  $K^0 = P^0$  as defined in (2.5). As a side remark we note that a similar argument can be used to show that it is immaterial whether extremal internal sectors are assigned the number  $\frac{1}{2}$  or  $N + \frac{1}{2}$ .

#### 4. Renormalization and the Existence Problem

As yet, the expressions derived in the preceding section have only a formal meaning. There remains the question of the existence of the integrals symbolized by the graphs.

The UV problem is concerned with the behaviour of the integrands at infinity in momentum space. It can be handled exactly as was done in V for the vacuum case. We note first that loop integrals over loops extending over more than one sector are finite because of the strong decrease of  $\tilde{D}_+(k)$  for  $k_0 \rightarrow -\infty$  and momentum conservation: primitively divergent subgraphs exist only within sectors, where they can

be treated with conventional methods. We choose BPHZ renormalization (see [7]), which introduces suitable subtractions in the integrand, before integrating, thereby avoiding the need for regularization. The subtractions are found by expanding the integrands of potentially dangerous subgraphs, the “renormalization parts,” into power series of sufficiently high degree in their external momenta. The temperature dependent part  $\tilde{C}$  of the propagator  $\tilde{D}_F$  decreases exponentially fast at large momenta and is innocent of any UV problems. We can therefore define the mentioned subtraction terms using  $\tilde{A}_F$  instead of  $\tilde{D}_F$  inside the renormalization parts, without destroying the UV convergence achieved by the subtraction. We also define the finite renormalizations of the BPHZ prescription to be  $T$ -independent, giving them the values needed to satisfy the normalization conditions at  $T = 0$  stated in V. These rules are the expression of the required  $T$ -independence of the renormalization prescription  $N$  in the field equation (1.1). An important effect of this limited subtraction is the emergence for  $T > 0$ , of lines connecting a vertex to itself. But they carry the propagator  $\tilde{C}(k)$ , hence the loop integral over  $k$  exists, and is the same in  $T^+$  and in  $T^-$  sectors.

But the UV problem is not the only existence problem we are faced with. There is also the problem of the local singularities of the integrand, again in  $p$ -space. The integrand is a product of distributions in variables which are quadratic functions of a complete set of independent external and internal momenta. The variables in a given subset of propagators may be dependent (i.e. their gradients in momentum space may be linearly dependent) on certain manifolds, in which case the product of propagators does not necessarily define a distribution. It is no longer clearly integrable even in the sense of distributions. Contrary to assertions found in the literature, the problem is not restricted to the case of two propagators separated by a self energy insertion, and thus both depending on the same variable. As an example of a more complex situation, consider the integral  $\int dk \tilde{D} \cdot (k) \tilde{D} \cdot (p - k)$  over a two-line loop, the dots standing for either  $+$  or  $F$ . The mass shell  $\delta$ 's in the two factors coalesce at  $p = 0$ , hence the integral diverges at that point, and this singularity in the external variable  $p$  (external to the considered 2-line subgraph) is not removed by renormalization. Closer inspection shows that the singularity is of first order. A chain of  $n$  two-line bubbles will then produce a singularity of order  $|p|^{-n}$ , which is not integrable for  $n \geq 4$ . Nor is it defined in another way as a distribution. The remaining integration over  $p$ , which may be an internal or external variable of the full graph, is therefore not defined: individual graphs containing such chains diverge. These divergences may cancel between graphs with the same scaffolding but different sector assignments of the vertices in the chain. But we are not aware of a proof to this effect, even for this still relatively simple example.

How is this problem solved in the vacuum case? This is easiest for the fully time ordered function  $\tilde{\tau}(P) = \tilde{\mathcal{W}}(P, +)$ . If  $m > 0$  this function is everywhere the boundary value of an analytic function in complexified variables  $p_i$ . The same is true for the Feynman propagator  $(k^2 - m^2 + i\varepsilon)^{-1}$ . For variables  $\{p_i\}$  in the domain of analyticity of  $\tilde{\tau}$  we can deform the integration contours for the internal variables  $k_j$  into the complex in such a way that they never meet the singularities at  $k_j^2 = m^2$ . The integrand is then a smooth function, and there are no problems with the local existence of the integral. But the presence of a  $\delta$ -term in  $\tilde{D}_F$  destroys analyticity, so that the method does not work for positive  $T$ . For  $m = 0$ ,  $\tilde{A}_F(k)$  is singular at  $k = 0$  and is there not a boundary value of an analytic function. The same holds for  $\tilde{\tau}$  at points where a partial sum of  $p_i$ 's vanishes. These singularities can lead to infrared divergences, which need special attention. An  $x$ -space method for achieving

this, and also for proving existence of the general  $\mathcal{W}_\sigma$ , has been described in V. It relies heavily on the fact that  $\Delta_+(\xi)$  is analytic in  $\text{Im}\xi^0 < 0$  and decreases at least of second order for  $\text{Im}\xi^0 \rightarrow -\infty$ . But  $D_+(\xi)$  is only analytic in the strip  $-\beta < \text{Im}\xi^0 < 0$ , so that this method can also not be extended to positive  $T$ .

At present the problem of existence of  $\mathcal{W}_\sigma$  remains unsolved.

### 5. Uniqueness

It will be shown that the expansion described in Sects. 3 and 4, assuming its existence in finite orders, is the only expression satisfying all the requirements.

Assume that uniqueness has been established up to order  $\sigma - 1$ . Let  $W_\sigma^1, W_\sigma^2$ , be two solutions of Eqs. (3.2) plus subsidiary conditions. Then their difference

$$h_\sigma(x_1, \dots, x_n) = W_\sigma^1(x_1, \dots) - W_\sigma^2(x_1, \dots) \tag{5.1}$$

satisfies the homogeneous equations

$$(\square_i + m^2)h_\sigma(\dots, x_i, \dots) = 0 \tag{5.2}$$

and all the subsidiary conditions. We want to derive the most general form of  $h_\sigma$ .

The Fourier transform  $\tilde{h}_\sigma(p_1, \dots, p_n)$  must contain a factor  $\delta(p_i^2 - m^2)$  for each  $i$ .

Define

$$h_\sigma(\dots [x, y] \dots) = h_\sigma(\dots, x, y \dots) - h_\sigma(\dots y, x \dots) \tag{5.3}$$

and let  $\tilde{h}_\sigma(\dots [p, q] \dots)$  be its Fourier transform. As in V we can prove that  $\tilde{h}_\sigma(\dots [p, q] \dots)$  has its support contained in the manifold  $p + q = 0$ , and must therefore be of the form

$$\hat{h}_\sigma(\dots [p, q] \dots) = \delta^4(p + q)\delta(p^2 - m^2)\{f_+(\mathbf{p}) + \varepsilon(p_0)f_-(\mathbf{p})\}, \tag{5.4}$$

where the dependence on the variables indicated by dots has been suppressed. By locality the support of the function  $\theta(x^0 - y^0)h_\sigma(\dots [x, y] \dots)$  is contained in the set  $(x - y)\varepsilon\bar{V}_+$ , the closed forward cone. Its Fourier transform

$$\frac{i}{2\pi}\delta^4(p + q)\left[\frac{f_+(\mathbf{p})}{\omega(\mathbf{p})}\frac{p_0}{(p_0 + i\varepsilon)^2 - \mathbf{p}^2 - m^2} + f_-(\mathbf{p})\frac{1}{(p_0 + i\varepsilon)^2 - \mathbf{p}^2 - m^2}\right],$$

considered as a function of  $p + q$  and  $p$ , is in  $p$  analytic in the forward tube  $\{\text{Im} p \in V_+\}$  and is polynomially bounded in the slightly smaller tube  $\{\text{Im} p - a \in V_+\}$  for any  $a \in V_+$  [10]. The  $p_0$ -dependent factors in the above expression have the correct analyticity. In order for this to be the case for the full expression, the  $\mathbf{p}$ -dependent terms  $\omega^{-1}f_+$  and  $f_-$  must be entire functions of  $\mathbf{p}$ . Furthermore, they must be polynomially bounded, hence they are polynomials. As a result we find

$$h_\sigma(\dots [p, q] \dots) = \delta^4(p + q)\delta(p^2 - m^2)\{\omega(\mathbf{p})F(\mathbf{p}) + \varepsilon(p_0)G(\mathbf{p})\} \tag{5.5}$$

with  $F$  and  $G$  polynomials in  $\mathbf{p}$ .  $F$  and  $G$  may also depend on the variables not shown explicitly, and factors only depending on those other variables have been omitted.

Consider next the double-commutator function

$$h_\sigma(\dots [[x, y], z] \dots) = h_\sigma(\dots [x, y]z \dots) - h_\sigma(\dots z[x, y] \dots). \tag{5.6}$$

By locality its support is contained in the set  $S = \{(x - z)^2 \geq 0 \text{ or } (y - z)^2 \geq 0\}$ . But by transforming (5.5) into  $x$ -space we find that the right-hand side of (5.6) depends on  $x$  and  $y$  only in the combination  $x - y$ , i.e. it is invariant under simultaneous translation of  $x$  and  $y$  by the same 4-vector. But no subset of  $S$ , apart from the empty set, is invariant under such a simultaneous translation. Therefore we find

$$h_\sigma(\dots [[x, y]z] \dots) = 0. \tag{5.7}$$

Finally we find from the KMS condition (2.6) that

$$\begin{aligned} (1 - \exp(-\beta p_1^0)) \tilde{h}_\sigma(p_1, \dots, p_n) &= \tilde{h}_\sigma(p_1, \dots, p_n) - \tilde{h}_\sigma(p_2, \dots, p_n, p_1) \\ &= \sum_{i=2}^n \tilde{h}_\sigma(p_2, \dots, [p_1, p_i] \dots, p_n), \end{aligned} \tag{5.8}$$

from which  $\tilde{h}_\sigma$  can be obtained through division by  $(1 - \exp(-\beta p_1^0))$ . Reinserting the resulting expression in the terms on the right-hand side of Eq. (5.8), and iterating this procedure sufficiently many times, we find that  $\tilde{h}_\sigma(p_1, \dots, p_{2n})$  must be of the form

$$\begin{aligned} &\tilde{h}_\sigma(p_1, \dots, p_{2n}) \\ &= \sum F(\mathbf{p}_{i_1}, \dots, \mathbf{p}_{i_n}) \prod_{\alpha=1}^n [1 - \exp(-\beta p_{i_\alpha}^0)]^{-1} \delta(p_{i_\alpha}^2 - m^2) \delta^4(p_{i_\alpha} + p_{j_\alpha}), \end{aligned} \tag{5.9}$$

where the sum extends over all partitions of the  $2n$  variables into  $n$  ordered pairs  $(p_{i_\alpha}, p_{j_\alpha})$ ,  $i_\alpha < j_\alpha$ , and  $F$  is a polynomial in  $\mathbf{p}_i$  and  $\omega(\mathbf{p}_i)$ . The odd-point functions vanish in  $\Phi^4$  theories.

Now, if  $F$  is a genuine polynomial, not a constant, then the expression (5.9) shows a bad high-energy behaviour for  $p_{i_\alpha}^0 \rightarrow \infty$ , which would destroy renormalizability. Hence  $F$  must be a constant, and  $h_\sigma$  has in  $x$ -space the form

$$h_\sigma(x_1, \dots, x_{2n}) = c_{\sigma n} \sum_{\text{pairings } \alpha} \prod D_+(x_{i_\alpha} - x_{j_\alpha}). \tag{5.10}$$

For dimensional reasons  $c_{\sigma n}$  must contain a factor  $\hbar^n$ , independently of  $\sigma$ . Insertion of such expressions in the right-hand side of the equation of motion (3.2) leads to terms of higher order in  $\hbar$ . This means that the homogeneous terms (5.10) are the contributions of lowest order to the  $2n$ -point functions. Since their functional form does not depend on  $\sigma$ , their contributions to  $W(x_1, \dots, x_{2n})$  can be summed over  $\sigma$ , in the sense of formal power series, to yield

$$c_n(g) \sum_x \prod D_+(\dots) \tag{5.11}$$

as term of lowest order in  $\hbar$ . The normalization condition for the 2-point function determines the value of  $c_1$  for  $T = 0$  to be  $c_1 = -i$ . With a different choice of  $c_1$  for  $T > 0$  our  $T$ -independent definition of the normal product  $N$  would no longer

remove the ultraviolet divergences, i.e. such a choice would destroy the validity of the field equation (1.1). Hence we must have

$$c_1 = -i \quad (5.12)$$

identically in  $g$  and  $T$ .

By our assumptions the expression (5.11) must satisfy the cluster property (2.3). Using induction with respect to  $n$  one finds the unambiguous result

$$c_n = (-i)^n. \quad (5.13)$$

This value does not depend on  $g$ , which proves the desired unicity

$$h_\sigma = 0 \quad (5.14)$$

for  $\sigma > 0$ . In the lowest order  $\sigma = 0$ ,  $h_0$  is the only nonvanishing contribution to the  $2n$ -point function. It has the conventional free-field form, as has already been anticipated in Sect. 3.

It is a curious aspect of this argumentation that it does not work in the free case  $g = 0$ , where renormalizability is not at stake. For the free theory one could arrive at the same results faster by postulating the usual canonical commutation relations. For interacting fields this method is rather less convincing, because in that case the CCR's have no rigorous meaning, since interacting fields cannot be restricted to sharp times.

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