

**FROM PERTURBATIVE  
TO CONSTRUCTIVE RENORMALIZATION**

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**PART I: INTRODUCTION****TO EUCLIDEAN FIELD THEORY**

## I.1 The ultraviolet problem

Quantum field theory is an attempt to describe the properties of elementary “point-like” particles in terms of relativistic quantum fields. It is now widely believed to offer a coherent mathematical framework for relativistic models (like the “standard  $U(1) \times SU(2) \times SU(3)$  model”). These models include all the particles and interactions observed up to now except gravity. Therefore, together with general relativity, field theory is the backbone of our current understanding of the physical world. In the future a new, more unifying framework may be adopted, like the currently promising superstring theory, which is a relativistic and quantum modelization of extended one dimensional objects instead of point-like particles; nevertheless even in this case it is extremely likely that field theory will remain important in many situations, just as classical mechanics is still today.

This situation is relatively recent. Until the 70’s the very statement that quantum field theory might provide a coherent mathematical framework at all was not widely accepted. The main doubts on the mathematical consistency of quantum field theory were due to the persistence of ultraviolet problems (and to the lack of successful models for strong interactions: QCD, the present field theory of strong interactions, did not exist). Let us sketch what these ultraviolet problems are and why they are important.

An ultraviolet problem is one which is due to the existence of arbitrarily small length scales, or equivalently of arbitrarily large frequencies in the Fourier analysis of a theory. Such problems are inherent to the formalism of quantum field theory, because it is a crucial assumption that the fields live on a continuous space time. One might wonder whether this continuity condition has anything to do with physics and whether the whole problem is not a mathematical artefact. After all it is reasonable to expect that space time will conserve its smoothness only until the Planck scale, where quantum aspects of gravitation might distort it significantly. This Planck scale might provide a physical ultraviolet cutoff; this is what seems to occur in superstring theories where it is conjectured that at least in perturbation, there are no ultraviolet divergences. However the Planck scale is much higher than the typical scales that field theory tries to describe, and is completely inaccessible to direct experiments.

In fact the most compelling reason for which we are interested in the continuous formulation of field theory is the same for which we are interested in the thermodynamic limit of statistical systems. In statistical mechanics this limit corresponds to systems of infinite volume. We know that in nature macroscopic systems are in fact finite, not infinite, but they are huge with respect to the atomic scale. The thermodynamic limit is an adequate simplification in this case, since it allows to give a precise mathematical content to the physically relevant questions (like dependence of the limit on boundary conditions, existence of phase transitions etc...). Since a limit has been taken, the power of classical analysis may be applied to these questions. It would be much harder and less natural to try to define the analogous notions for a large finite system, just as it is difficult and often inappropriate to make discrete approximations to some typically continuous mathematics like topology.

From this point of view the ultraviolet problem appears central and inescapable

in field theory; a limit has to be performed, whose existence is critical for interest in the corresponding mathematical formalism.

Historically quantum field theory was plagued by two successive ultraviolet “diseases” which raised doubts on the existence or consistency of the ultraviolet limit. In both cases the situation looked bad for many years until a way out of the crisis was found. The first and most famous ultraviolet disease has been recognized almost since the birth of quantum field theory. It is the occurrence of divergences due to the integration over high momenta in the loops of Feynman integrals. In the  $\phi_4^4$  theory which will be discussed soon, one of the simplest of these divergences is the divergence of the second order graph which we call the “bubble” (see Fig. I.1.1). By momentum conservation the amplitude for this graph is only a function of  $k = k_1 + k_2$ . In Euclidean space, this amplitude (apart from combinatoric coefficients discussed later) is given by the integral:

$$\int \frac{d^4 p}{(p^2 + m^2)[(p + k)^2 + m^2]} \quad (\text{I.1.1})$$

which diverges logarithmically for large values of  $p$ . (Similar divergences occur of course in the more physical theory of quantum electrodynamics).

Around 1950, this disease was cured by the invention of perturbative renormalization by Feynman, Schwinger, Tomonaga, Dyson and others (see e.g. [Dy1]). Basically it amounts to a redefinition of the physically observable parameters of the theory which pushes the infinities into unobservable “bare” parameters. It took more than a decade to put this perturbative theory of renormalization on a completely firm mathematical basis. Roughly speaking, the main result states that theories which are renormalizable from the naive power counting point of view can indeed be renormalized without changing the formal structure of the Lagrangian. More precisely one can replace the bare parameters of the Lagrangian by formal power series in the renormalized parameters (usually the coupling constant), so that the resulting perturbative expansion in the renormalized coupling is finite to all orders, as a formal power series. We call this important theorem the BPH theorem (Bogoliubov- Parasiuk-Hepp [BP][He1]) although as usually it incorporates a lot of former work and was followed by important extensions and refinements. It was somehow a surprise to discover that this theorem, developed for quantum electrodynamics or the  $\phi_4^4$  model, remained also true for non abelian gauge theories [tH1][LZ], in which case it is highly non trivial to check that the counterterms required do not break gauge invariance, i. e. can be absorbed in a redefinition of the field strength and of the coupling constant. By power counting analysis alone, this would not be true: one has to incorporate additional information coming from Slavnov identities or BRS invariance [Tay][Sla][BRS]. The traditional proof (see [IZ]) relies on a dimensional regularization in which gauge invariance is maintained, but this approach has some drawbacks: the dimensional regularization is complicated and cannot be used up to now in a constructive (non perturbative) program.

But quantum electrodynamics, the only firmly established field theory has been plagued by another ultraviolet problem, raised in particular by L. Landau and other physicists of the Russian school. We call it the “renormalon” problem, although this name was introduced much later. It does not occur any more at

the level of individual Feynman graphs, but it affects the perturbative series as a whole. In  $\phi_4^4$  (which is similar to electrodynamics in this respect), there are several ways to discover the problem. One of them is to consider the leading-log behavior of the renormalized 4 point function at  $n$ -th order,  $S_n^{4,R}$ , and large external momenta. Let  $m$  be the mass of the particles. One finds (in Euclidean space):

$$S_n^{4,R}(k) \simeq_{k \rightarrow \infty} g_R^n [\beta_2 \log(k/m)]^{n-1} \quad (\text{I.1.2})$$

where  $g_R$  is the renormalized coupling constant and  $\beta_2$  is a numerical coefficient (for single component  $\phi_4^4$  it is  $9/2\pi^2$  or  $3/16\pi^2$  if one writes, as usual, the interaction as  $\phi^4/4!$ ). Various ways of playing with formula (I.1.2) give rise to various troubles, all related. When the 4 point function is inserted into a convergent loop like the triangular 6-point graph of Fig. I.1.2, at 0 external momenta, one obtains contributions to the  $n$ -th order of perturbation theory for the 6 point function proportional to

$$g^n \int_{|k| \geq m} \frac{d^4 k}{(k^2 + m^2)^3} [\beta_2 \log(k/m)]^{n-3} \simeq (n-3)! g^n \beta_2^{n-3} \quad (\text{I.1.3})$$

These contributions are not summable over  $n$  (since they add up with the same sign they are also not Borel summable). Therefore the renormalon problem appears as a difficulty to sum up perturbation theory. However one might also consider the asymptotic behavior in  $k$  of the bare (unrenormalized) 4-point function:

$$S_{n,k}^{4,B}(O) \simeq_{k \rightarrow \infty} (-g_B)^n [\beta_2 \log(k/m)]^{n-1} \quad (\text{I.1.4})$$

where  $g_B$  is the bare coupling constant. Since in the theory of perturbative renormalization this function should be the counterterm of the theory, i.e. the difference between the renormalized coupling  $g_R$  and the bare coupling  $g_B$ , one gets the formula:

$$g_R = g_B - \sum_{n=2}^{\infty} (-g_B)^n [\beta_2 \log(k/m)]^{n-1} = \frac{g_B}{1 + g_B \beta_2 \log(k/m)} \quad (\text{I.1.5})$$

Therefore if  $g_B > 0$ ,  $\lim_{k \rightarrow \infty} g_R = 0$  (no matter how  $g_B$  is chosen as a function of  $k$ ). This is the phenomenon of “charge screening” called also the triviality problem for  $\phi_4^4$  or QED. In the ultraviolet limit, only the trivial noninteracting theory seems to exist. Still another possibility is to invert formula (I.1.5) to obtain

$$g_B(k) \simeq \frac{g_R}{1 - g_R \beta_2 \log(k/m)}. \quad (\text{I.1.6})$$

Keeping  $g_R$  fixed,  $g_B$  becomes negative at an energy of order  $me^{+\frac{1}{\beta_2 g_R}}$ , therefore the theory should become unstable or inconsistent at this “Landau” energy scale. The singularity in (I.1.6) is sometimes called a Landau ghost. But of course one should not really trust (I.1.6). One should rather remark that as  $k \rightarrow \infty$ ,  $g_B$  increases and becomes of order unity around the Landau energy, after which perturbation theory itself and in particular (I.1.2)-(I.1.6) should no more be valid. The behavior of  $g_B$  is turned into a strong coupling problem, intractable up to now, except perhaps by Monte Carlo simulations.

Let us summarize the characteristic features of the renormalon problem. As remarked already, it affects the summability of perturbation theory as a whole (therefore physicists might call it a non-perturbative problem, although we would prefer to consider it as a “strong coupling” problem). It is truly an ultraviolet disease, because it does not occur in the theory with fixed ultraviolet cutoff, no matter how large. It is easy for instance to check that the  $n!$  behavior of (I.1.3) does neither appear in the bare nor in the renormalized perturbative series with fixed ultraviolet cutoff. Finally the name “renormalon problem” although perhaps awkward, is justified because the disease arises from the introduction of counterterms, hence from the use of perturbative renormalization. As will indeed be discussed at length in this book, the perturbative theory of renormalization cures the first problem of infinities too well. It introduces both some pieces of counterterms which we call “useful” because they make the renormalized Feynman amplitudes finite, but also some pieces of counterterms which we call “useless”. These useless counterterms do not cure any divergence. Furthermore they are the ones which are responsible for the renormalon problem!

To observe the distinction between useful and useless counterterms, one needs some detailed Fourier analysis such as the one provided by phase space decomposition. Then it appears that the sole reason for introducing useless counterterms is the locality requirements and the insistence upon writing the renormalized series is written in terms of fixed renormalized constants which do not depend on the energy scale. In this book we will argue that one should drop this restriction, and adopt a more effective perturbation theory with an infinite number of scale dependent “effective constants”. In this way we rediscover that the renormalization group point of view is the right way to investigate the renormalon problem.

Historically the renormalon disease and its investigation by the invention of the renormalization group was not discussed exactly in these terms; the emphasis was on the invariance of the theory under changes of the (arbitrary) subtraction scale (this invariance lead to the very name “renormalization group”) and on high energy asymptotic behavior, not on  $n!$  behavior at large order and the problem of summing up perturbation theory. Of course both points of view are closely related. Anyway the problem was serious enough to raise again doubts on the ability of quantum field theory to be consistent through the 60’s.

As a reaction the constructive program was launched, and in the early 70’s a major milestone was reached when superrenormalizable theories of various type were built and checked to be free of inconsistencies. Reviews or books on this first period of constructive theory are [Er1][Si1][GJ2]. Although the results and techniques used have proved very influential in many areas of physics and mathematics, this success was nevertheless not of specific relevance to the renormalon problem. For theoretical physicists, a convincing way to escape this problem was really found with the major discovery of asymptotic freedom in nonabelian gauge theories [Po][GW]. This occurred just at the right time to complete the spectacular rebirth of quantum field theory: non abelian gauge theories had been developed as realistic models for the electroweak interaction and had been shown to be perturbatively renormalizable. It gave in turn a major impetus to adopt them to describe strong interactions as well.

Asymptotic freedom occurs when the coefficient  $\beta_2$  in (I.1.2) is negative. As

a result equations (I.1.5)-(I.1.6) are inverted. It is now the bare charge which is screened. At large energies the particles behave like free point-like objects (hence the name of asymptotic freedom). Looking at (I.1.3) and changing the sign of  $\beta_2$  we see that the renormalon problem still prevents “ordinary” summation, but the corresponding contributions now become alternate; therefore an other type of summability, like Borel summability, becomes possible.

The discovery of asymptotic freedom in non-abelian gauge theories convinced the theoretical physics community that these quantum field theories are indeed mathematically consistent. Many physicists believe that there is no longer any surprises to arise on the ultraviolet problem for gauge theories (see however [tH4] for an exception). But this belief has yet to be better substantiated by a non perturbative, mathematically rigorous analysis.

To understand rigorously the concept of asymptotic freedom beyond perturbation theory, one way is to construct first some consistent models of renormalizable theories with such a behavior. There exist models of this kind simpler than the non abelian gauge theories in 4 dimensions, namely fermionic models in two dimensions with many components and a quartic interaction [MW] [GrNe]. Also the  $\phi_4^4$  model, although not asymptotically free in the ultraviolet direction, is asymptotically free in the infrared direction. Although the corresponding constructive problem is a problem of statistical mechanics rather than field theory (the ultraviolet cutoff is not removed), it is very similar in mathematical structure. This road has been followed by K. Gawedzki and A. Kupiainen [GK2-3-4] and by J. Feldman, J. Magnen, R. Sénéor and the author [FMRS3-4-5], who, with somewhat different technical tools, succeeded in building these models. In fact it is the main goal of this book to present in a more systematic and accessible form than the original papers the technique of multiscale or phase space expansion, as developped and applied in our collaboration with J. Feldman, J. Magnen and R. Sénéor. This technique originated in the constructive work of Glimm and Jaffe [GJ1].

To extend our rigorous understanding of asymptotic freedom to non abelian gauge theories in a finite volume is the next natural challenge. (The large volume problem is indeed a different one, where one has to deal with a strong coupling problem and physical issues like quark confinement which seem still much farther from a rigorous analysis). This ultraviolet consistency of non-abelian gauge theories on compact manifolds is not only a key issue in theoretical physics, but is also becoming one in geometry, in particular since E. Witten related Donaldson’s invariants to (still formal) functional integrals of some (supersymmetric) Yang-Mills theory [Wit]. Various versions of Yang-Mills theories now seem to be the link between the most fascinating problems of geometry (homotopy and knots in three dimensions, symplectic geometry and conformal theories in two dimensions and differential geometry in four dimensions). A constructive understanding of the corresponding functional integrals would be therefore a major progress in pure mathematics.

The problem has been attacked first by T. Balaban [Ba2-9]. In a very impressive sequence of papers, completed recently, he establishes an ultraviolet stability bound for the effective action of a lattice gauge theory after iterating a large number of clever block-spin transformations. From this result it is expected that the continuum limit of gauge invariant observables like Wilson loops can be con-

structed. Hence at least these observables should be free from inconsistencies and this is a very important result. This work is now followed by a related program of P. Federbush [Fe2-7][FedW], still in progress, which uses also a lattice regularization and phase space cells.

But even after completion of these works, many questions remain in particular because the expectation values of the fields of the theory (in a particular gauge) are not built in the lattice-based approach. It is important to know whether such expectation values, which are not gauge invariant, can also be built or, if they cannot, to understand better why. The main difficulty for this other program lies in a lack of positivity of the gauge-fixed functional integrals, related to the Gribov problem [Gri]. Progress in understanding this difficult problem has been slow up to now, but interesting as well as surprising results may lie ahead [Zw1-3][DeZw].

After this quick historic overview of the ultraviolet problem in field theory, let us describe the structure of this book.

The first part is devoted to some introductory material on field theory, on the  $\phi^4$  model which is the training ground for most of the book, and on perturbation theory and Feynman graphs. We keep this part very brief; although we try to be reasonably self-contained, we assume indeed some familiarity of the reader with quantum field theory, for instance with [IZ].

Then in the second part we apply the idea of “phase space chopping” to the study of perturbation theory, focusing on renormalization and using the  $\phi_4^4$  model for simplicity. We derive a uniform BPH theorem which adds reasonable estimates to the finiteness content of the original BPH theorem; the first theorem of this kind appeared in [dCR1]. Then we show that the analysis of the bounds obtained leads naturally to reshuffle the bare or renormalized perturbation theory into a better form, the effective perturbation theory. In the case of simplified asymptotically free models for which the number of graphs at ordre  $n$  is not too large, like the wrong sign planar  $\phi_4^4$  model [tH5-7], [Ri1], this reshuffling is even sufficient to construct the model. This part is closed by a discussion of the large order *behavior* of perturbation theory, a problem truly at the border between perturbative and constructive techniques.

In the third part, devoted to constructive theory, we start with an introduction to the key techniques of cluster and Mayer expansion; these techniques extend in a natural way to phase space under the name of “multiscale expansions”. Then we describe how to apply them first to the infrared limit of critical  $\phi_4^4$ , for which we show how positivity of the interaction can be used to handle the so called “domination” or large field problem, and then to the construction of the “Gross Neveu” model in two dimensions, which is a genuine renormalizable field theory. The construction of the non renormalizable Gross-Neveu model in three dimensions, which requires interesting additional techniques, is also sketched.

We conclude with a description of a tentative approach to the problem considered above, namely the construction of the continuum limit for the gauge-fixed non-abelian gauge theories  $YM_4$  in a finite (small) volume. We restrict for simplicity to the pure  $SU(2)$  theory. We want to construct the continuum limit of the gauge fixed theory in a regular gauge, because this theory has a regular propagator which is compatible with phase space chopping. Although the phase space chopping does not respect gauge invariance, it should nevertheless be possible to

recover gauge invariance in the limit, at least in the form of Slavnov identities which express invariance under gauge transformations continuously connected to the identity. This is because only the relevant or marginal effects of gauge breaking cutoffs matter for the final theory, and these effects can be compensated by means of a finite number of appropriate non gauge invariant counterterms. There are some stability requirements for these counterterms, which as we show can be met. This approach has lead only to modest results up to now, but at least as a by-product it gives a control of perturbative renormalization of non-abelian gauge theories with reasonable bounds, which avoids the traditional use of dimensional regularization, which is not adapted to constructive purposes.

The difficulties that we met in our program may indicate that the ultraviolet behavior of the non-abelian gauge theories, at least at the level of the vacuum expectation values in a regular gauge, is *not* the one expected from the perturbative renormalization group predictions, as is advocated in [DeZw]. This important point clearly calls for a rigorous clarification.

To conclude, let us apologize sincerely to the many experts in all these areas whose work is not properly cited or accounted for in this book. Neither on perturbative nor on constructive theory, neither at the level of the subjects treated nor at the level of the references does this book intends to be an extensive or even fair review. The main reason is that we do not feel able to report properly on something else than the particular techniques and point of view that we have personally used, which we know may not be always the most elegant or the most powerful ones. As a result, this book may be considered more as a “guided tour” of the author’s favourite subjects than as an exhaustive review on perturbative or constructive field theory.

Nevertheless we want to mention at this stage that concerning the part of this book on perturbative renormalization, an other elegant formalism using phase space chopping was developed in parallel to ours by Gallavotti and Nicoló, and later by Feldman, Hurd, Rosen and Wright. This formalism is very closely related to the one presented here, the main difference being that it avoids the use of Feynman graphs, making in a sense the combinatoric of cancellations more transparent. Here we stick to the point of view of Feynman graphs, which we think are presumably familiar to the potential reader, and refer for this other approach to [Gal][GaNi][FHRW].

Concerning the third part of this book, we mentioned already that a renormalization group solution to many of the same problems was worked out simultaneously by K. Gawedzki and A. Kupiainen [GK2-4], using the formalism of block spin transformations. These beautiful works and many others related will not be described at all, again for lack of competence but not of interest!

A last word is in order on the mathematical level of rigor. We intend to describe only rigorous results, but we do not provide always complete proofs of all the statements. Also we prefer always the particular to the general, for pedagogical reasons. We try to convey better the underlying ideas by phrasing as much as possible our arguments in ordinary words rather than in fully formalized equations. We know that this is dangerous and that the ratio of equations to words in scientific papers has been sometimes compared to that of signal to noise. In any case we did honestly our best to mention all the difficulties, even when they are apparently

purely technical. In conclusion we hope that this book can be useful, in particular to the beginner, by explaining the natural link between perturbation theory, which is conceptually simpler, and constructive field theory, which had for too long the reputation of being a much harder subject.

## I.2 Euclidean field theory. The O. S. axioms

From the classical work of the founding fathers of axiomatic field theory, we learn that the minimal mathematical requirements or axioms for a field theory are conveniently expressed in terms of the vacuum expectation values of products of the field operators, the “Wightman functions” [SW]. From these quantities one could also, at least in principle, compute more physical quantities like the  $S$  matrix. We could start directly from the Wightman axioms or their Euclidean counterpart, the Osterwalder-Schrader axioms, but we prefer to motivate them first with a brief sketch of their relation to the  $S$  matrix formalism, without any attempt to mathematical rigor. For this sketch, we follow [dC]; for a more complete study on general aspects of quantum field theory, we refer to [BS],[BD] and [IZ].

In ordinary (Minkowski) quantum field theory, the fields are technically “operator valued distributions”, i.e. they take their values in some set of operators on the Hilbert space of physical states,  $H$ . The main physical properties that are required for field theory are relativistic covariance and microcausality. Therefore  $H$  should bear a unitary representation of the Poincaré group, for which the only invariant state is the vacuum. The generators of the Poincaré group, the momentum operators  $P$  should have positive norm:  $P^2 \geq 0$ . And since signals do not propagate faster than light, operators which are smeared with test functions whose supports are “space like” separated should commute (or anticommute for fermions).

Since field theory is a second quantized formalism, the Hilbert space is a Fock space:

$$H = \bigoplus_{n=0}^{\infty} H_n \quad (\text{I.2.1})$$

$H_0$ , of dimension 1, is generated by a particular vector  $\psi_0$ , called the vacuum.  $H_1$  is the space of 1-particle states, and  $H_n$ , the space of n-particle states, is the n-th tensor product of  $H_1$ , symmetrized for bosons and antisymmetrized for fermions. For instance for free massive scalar bosons,  $H_1$  is generated by a complete orthonormal set  $f_i$  of positive energy solutions of the Klein-Gordon equation:

$$(\partial_\mu \partial^\mu + m^2)f = 0 \quad (\text{I.2.2})$$

The scalar product in this case is

$$\delta_{ij} = \langle f_i, f_j \rangle = i \int d^3x f_i^* \bar{\partial}_0 f_j = \int d\mu_p \tilde{f}_i^* \tilde{f}_j \quad (\text{I.2.3})$$

where  $f_i = \int d\mu_p \tilde{f}_i(p) e^{-ip \cdot x}$ ;  $d\mu_p$  is the Lorentz invariant measure:

$$d\mu_p = \frac{d^3p}{(2\pi)^3 2p_0} = \frac{d^4p}{(2\pi)^4} 2\pi \delta(p^2 - m^2) \theta(p_0) \quad (\text{I.2.4})$$

and we use the convention  $u\bar{\partial}_0 v = u(\partial_0 v) - (\partial_0 u)v$ . The free scalar bosonic field  $\phi$  also satisfies the Klein Gordon equation (I.2.2), which derives from the free Lagrangian:

$$L_0(\phi) = \frac{1}{2} [\partial_\mu \phi \partial^\mu \phi - m^2 \phi^2] \quad (\text{I.2.5})$$

and it can be developed into creation and annihilation operators which satisfy canonical commutation relations:

$$\phi(x) = \int d\mu_p [a(p)e^{-ip \cdot x} + a^\dagger(p)e^{ip \cdot x}] \quad (\text{I.2.6})$$

$$[a(p), a(p')] = [a^\dagger(p), a^\dagger(p')] = 0; \quad [a(p), a^\dagger(p')] = 2p_0(2\pi)^3\delta_3(\vec{p} - \vec{p}') \quad (\text{I.2.7})$$

The action of the field on the Hilbert space is then best described by smearing these operators with solutions  $f$  of the Klein-Gordon equation. We can define

$$a_f = \int d\mu_p \tilde{f}(p)a(p), \quad a_f^\dagger = \int d\mu_p \tilde{f}(p)a^\dagger(p) \quad (\text{I.2.8})$$

The vacuum is then annihilated by any annihilation operator  $a_f$ , and a generating set for the full Hilbert space  $H$  is obtained by the action of a finite number of creation operators  $a_f^\dagger$  on the vacuum (which is therefore called a cyclic vector). The action of  $a_f$  and  $a_f^\dagger$  is then determined by the commutation relations coming from (I.2.7):

$$[a_f, a_{f'}] = [a_f^\dagger, a_{f'}^\dagger] = 0; \quad [a_f, a_{f'}^\dagger] = \langle f, f' \rangle \quad (\text{I.2.9})$$

To solve the Klein-Gordon equation one can introduce Green's functions (advanced, retarded or the symmetric one called Feynman's propagator), and the vacuum expectation value of a time ordered product of field operators (or generalized, N-point Green's functions) is given by a sum over pairings of these fields of the corresponding product of propagators. Such pairings are called Wick contractions. The result is called Wick's theorem. We will find this rule again in the context of gaussian integration.

The theory of the free field is therefore completely explicit. But we are in fact in search of an interacting quantum field theory. We look for field theories which admit a simple Lagrangian, polynomial in the fields and their derivatives. For scalar fields this Lagrangian will be decomposed into the free piece  $L_0$  given by (I.2.5) and a higher order polynomial  $L_i$  which is the interaction. Let us use  $\Phi$  for the interacting field, to distinguish it from the free field  $\phi$ . For instance in the  $\Phi^n$  theory the interaction is  $g\Phi^n$  and the corresponding field equation is a non linear generalization of the Klein-Gordon equation:

$$(\partial_\mu \partial^\mu + m^2)\Phi = ng\Phi^{n-1} \quad (\text{I.2.10})$$

The traditional approach to a collision process is to start from a system of free particles at time  $-\infty$  and to end up also with a system of free particles at time  $+\infty$ . The corresponding asymptotic spaces  $H_{in}$  and  $H_{out}$  should be therefore isomorphic to  $H$  and the collision process should be represented by a unitary matrix called the  $S$  matrix mapping  $H_{in}$  to  $H_{out}$ . Cross sections are then obtained directly from the matrix elements of  $S$  in a suitable basis. Unitarity of  $S$  is necessary so that probabilities of outgoing states add up to 1, no matter what the incoming state is. One requires also invariance of  $S$  under the Poincaré group, and stability of the vacuum (which should therefore be also invariant by  $S$ ). However many mathematical problems arise with this approach. First (I.2.10) contains a multiplication of distributions, which is generally an ill-defined operation. Secondly,

solutions of the non-linear equation (I.2.10) cannot be asymptotic to free fields hence to solutions of the linear one, except in a certain weak sense. The result of [LSZ] is that under such a suitable asymptotic condition, there are “reduction formulae” which express the matrix elements of  $S$  in terms of the Green’s functions  $G_N$  (or time ordered vacuum expectation values) of the interpolating (interacting) field  $\Phi$ :

$$G_N(z_1, \dots, z_N) = \langle \psi_0, T(\Phi(z_1), \dots, \Phi(z_N))\psi_0 \rangle . \quad (\text{I.2.11})$$

The Gell-Mann-Low formula gives in turn these functions as vacuum expectation values of a similar product of free fields with  $e^{iL_i}$  inserted:

$$G_N(z_1, \dots, z_N) = \frac{\langle \psi_0, T\left[\phi(z_1), \dots, \phi(z_N)e^{i \int dx L_i(\phi(x))}\right]\psi_0 \rangle}{\langle \psi_0, T(e^{i \int dx L_i(\phi(x))})\psi_0 \rangle} \quad (\text{I.2.12})$$

This formula is difficult to justify because the usual argument, based on the so called “interaction picture” is wrong: by a theorem of Haag, there is no way to relate the free field  $\phi$  to the interacting one  $\Phi$  by a unitary operator [SW]. Nevertheless the Gell-Mann-Low formula can be rigorously justified at least at the level of perturbation theory (in the sense of formal power series in the coupling  $g$  appearing in front of  $L_i$ ). Branching the coupling in an adiabatic way one finds (I.2.12) up to renormalization ambiguities, which mean that a certain finite part has to be taken in (I.2.12) [EG].

In the functional integral formalism proposed by Feynman [FH], the Gell-Mann-Low formula is itself replaced by a functional integral in terms of an (ill-defined) “integral over histories” which is formally the product of Lebesgue measures over all space time. It is interesting to notice that the integrand appearing in this formalism contains the full Lagrangian  $L = L_0 + L_i$ , not only the interacting one. The corresponding formula is the Feynman Kac formula:

$$G_N(z_1, \dots, z_N) = \frac{\int \prod_j \phi(z_j) e^{i \int L(\phi(x)) dx} D\phi}{\int e^{i \int L(\phi(x)) dx} D\phi} \quad (\text{I.2.13})$$

The functional integral has potentially many advantages. First it relates Wick theorem to the rules of gaussian integration, and makes therefore perturbation theory very transparent. The fact that the full Lagrangian appears in (I.2.13) is interesting when symmetries of the theory are present which are not separate symmetries of the free and interacting Lagrangians, as is the case for non-abelian gauge theories. It is also well adapted to constrained quantization, and to the study of non-perturbative effects. But for a long time nobody knew whether some day some rigorous meaning could be attached to it.

In fact the difficulty of making rigorous sense out of either the Gell-Mann-Low or the Feynman-Kac formula lead to the investigation of what minimal properties should be expected from the Green’s functions of a field theory, no matter from which Lagrangian or formula they may come from. Several axiomatic schemes were developed. The time ordered Green’s functions can also be computed from ordinary (not time ordered) vacuum expectation values of products of the fields operators, called Wightman functions. It is for these functions that the most popular

axiomatic scheme, the Wightman axioms was formulated. It includes a regularity assumption (temperedness), a relativistic transformation law, spectral conditions, hermiticity, local commutativity, a positiveness condition and a cluster property. The “Wightman reconstruction theorem” [Wig1] ensures that from such properties one can reconstruct in a unique way a Hilbert space with a unitary representation of the Poincaré group, positive spectrum condition for the momentum operators, the correct properties for the domain, regularity and transformation law under Poincaré of the field operators, local commutativity (or microscopic causality) and uniqueness and cyclicity of the vacuum. This Hilbert space formulation is called the Garding-Wightman axiomatic scheme. But in absence of interacting models, all these axioms could be checked only for free fields. In the late sixties, the constructive program was launched in order to provide at least some non trivial models for these axioms, and the simplest thing to do was to return to Lagrangians and to formulas (I.2.12) or (I.2.13).

There is a deep analogy between the Feynman Kac formula and the formula which expresses correlation functions in classical statistical mechanics. For instance, the correlation functions for a lattice Ising model are given by

$$\left\langle \prod_{i=1}^n \sigma_{x_i} \right\rangle = \frac{\sum_{\{\sigma_x=\pm 1\}} e^{L(\sigma)} \prod_i \sigma_{x_i}}{\sum_{\{\sigma_x=\pm 1\}} e^{L(\sigma)}} \quad (\text{I.2.14})$$

where the index  $x$  runs over the discrete points of the lattice, and  $L(\sigma)$  in the simplest case contains only nearest neighbor interactions and possibly a magnetic field  $h$ :

$$L(\sigma) = \sum_{|i-j|=1} J \sigma_i \sigma_j + \sum_i h \sigma_i \quad (\text{I.2.15})$$

By analytically continuing (I.2.13) in time to the Euclidean points, it is possible to complete the analogy with (I.2.14), hence to establish a firm contact with statistical mechanics. This idea also allows to give a rigorous meaning to the Euclidean path integral, at least for a free bosonic field. Indeed the corresponding Euclidean measure  $Z^{-1} e^{-\int L_0(\phi(x)) dx} D\phi$ , where  $Z$  is a normalization factor, can be defined easily as a gaussian measure on the Schwartz space  $S'$  of rapidly decreasing distributions, using the general theory of such measures and a theorem by Minlos [Si1]. This is simply because  $L_0$  is a quadratic form of positive type.

The Green’s functions continued to Euclidean points are called the Schwinger functions of the model, and are given by the Euclidean Feynman-Kac formula:

$$S_N(z_1, \dots, z_N) = Z^{-1} \int \prod_{j=1}^N \phi(z_j) e^{-\int L_i(\phi(x)) dx} d\mu_0(\phi) \quad (\text{I.2.16})$$

$$Z = \int e^{-\int L_i(\phi(x)) dx} d\mu_0(\phi) \quad (\text{I.2.17})$$

This formula is still formal, even in a finite volume; for instance the factor  $e^{-\int L_i(\phi(x))}$  might be ill-defined since it involves multiplying distributions, or might not decrease at large  $\phi$ , in the case of an “unstable” potential. However the gain from

an oscillating factor in (I.2.13) to a real one in (I.2.15), turns out to be enormous in practice to make sense out of such a formula. When the potential is stabilizing, we get direct decrease in modulus of the integrand in the functional integral, something much easier to use than a rapidly oscillating factor.

Euclidean formulation also helps to understand and control the classical limit: classical solutions are the fields for which the action is stationary. In the classical limit in Minkowski space, to show that the functional integral is dominated by the contributions near these classical solutions involves a stationary phase method; but in Euclidean space it is a steepest descent method, usually a much easier problem to control.

But Euclidean field theory would not have attracted so much interest, were it not for theorems which allow to go back from Euclidean to Minkowski space. There are several axiomatic schemes in Euclidean space which ensure such a recovery of a Wightman theory and the corresponding theorems are discussed quite extensively in [Sil]. We restrict us here to a brief survey of the Osterwalder-Schrader (O.S.) axioms [OS]. These axioms are expressed as properties of the Schwinger functions  $S_N$ , not of the field  $\phi$ . Therefore, strictly speaking, they are a kind of “Euclidean formulation of field theory” rather than “Euclidean field theory”. We will forget this subtlety. They are an efficient scheme for constructive purposes and in this book they will be particularly convenient since the Schwinger functions are the natural objects under study, both in the perturbative and constructive parts.

The O.S. axioms include five properties:

- OS1 A regularity property
- OS2 Euclidean covariance
- OS3 O.S. positivity
- OS4 Symmetry
- OS5 Cluster property

OS1 is purely technical; for the analytic continuation to Minkowski space to work one must check that the set of moments  $S_N$  does not grow too quickly with  $N$ . This axiom is usually very easy to check in constructive theory.

OS2 simply states that the Schwinger functions are invariant under a global Euclidean transformation, in the case of a scalar bosonic field. Of course for theories with spinors, the proper law of transformations of the spinor fields under the Euclidean group must be taken into account. After analytic continuation, this property ensures the proper covariance under the Poincaré group.

OS3 is the most interesting axiom. It states that the expectation value of a function (like a polynomial) of the fields multiplied by the same function after reflection on any hyperplane (and complex conjugation) must be positive. This key property ensures after continuation the locality axiom and positivity of the Hilbert space in the Wightman axiomatic scheme.

OS4 is just full symmetry for the Schwinger functions under permutations of the external arguments, as should be the case for bosons. For fermions, of course, one should replace this rule by antisymmetry.

OS5 states that the Schwinger functions asymptotically factorize when two sets of arguments are taken far apart. This ensures the unicity of the vacuum in the Wightman axioms. In theories where all particles are massive, the clustering

is exponential with the separation distance. For the two point function, the rate of decay is called the mass gap.

There are in fact several technical possibilities for OS1, depending on whether one requires full equivalence with the Wightman's axioms or simply a reconstruction theorem  $OS \rightarrow W$ . For a discussion of this and a full mathematical presentation of these axioms we refer to [OS],[Si1].

The important result is:

**Osterwalder Schrader reconstruction theorem**

Any set of functions satisfying OS1-OS5 determine a unique Wightman theory whose Schwinger functions they are.

From now on we forget the Minkowski space and all the background briefly reviewed above. We always assume that we are in a  $d$  dimensional Euclidean space  $\mathbb{R}^d$ . Our starting point is the Euclidean Feynman-Kac formula; our goal is to make rigorous sense out of it, and to check the validity of Osterwalder-Schrader axioms for the corresponding Schwinger functions.

### I.3 The $\phi^4$ Model

The simplest interacting field theory is the theory of a one component scalar bosonic field  $\phi$  with quartic interaction  $g\phi^4$  ( $\phi^3$  which is simpler looks unstable). In  $\mathbb{R}^d$  it is called the  $\phi_d^4$  model. For  $d = 1, 2, 3$  the model is superrenormalizable and has been built by constructive field theory. For  $d = 4$  it is renormalizable in perturbation theory. Although a constructive version may not exist [Aiz],[Frö], it remains a valuable tool at least for a pedagogical introduction to renormalization theory.

Formally the Schwinger functions of the  $\phi_d^4$  are the moments of the measure:

$$d\nu = \frac{1}{Z} e^{-(g/4!) \int \phi^4 - (m^2/2) \int \phi^2 - (a/2) \int (\partial_\mu \phi \partial^\mu \phi)} D\phi \quad (\text{I.3.1})$$

- $g$  is the coupling constant, usually assumed positive or complex with positive real part;
- $m$  is the mass; it fixes an energy scale for the theory;
- $a$  is the wave function constant. We often assume it to be 1.
- $Z$  is a normalization factor which makes (I.3.1) a probability measure.
- $D\phi$  is a formal product  $\prod_{x \in \mathbb{R}^d} d\phi(x)$  of Lebesgue measures at every point of  $\mathbb{R}^d$ .

Remark that the quadratic piece in the exponential of (I.3.1) describes the free propagation of particles with a positive definite propagator  $(p^2 + m^2)^{-1}$  in Fourier space. This suggests a first improvement of (I.3.1) towards mathematical respectability. Consider the translation invariant propagator  $C(x, y) \equiv C(x - y)$  (with slight abuse of notation), whose Fourier transform is

$$C(p) = \frac{1}{(2\pi)^d} \frac{1}{p^2 + m^2} \quad (\text{I.3.2})$$

We can use Minlos theorem and the general theory of gaussian processes [Er1][Si1] to define  $d\mu(\phi)$  as the gaussian measure on  $S'(\mathbb{R}^d)$  whose covariance is  $C$ . With this definition, the measure (I.3.1) should be equal to

$$\frac{1}{Z} e^{-(g/4!) \int \phi^4} d\mu(\phi) \quad (\text{I.3.3})$$

The Schwinger functions are then given by the still formal expression:

$$S_N(z_1, \dots, z_N) = \frac{1}{Z} \int \phi(z_1) \dots \phi(z_N) e^{-(g/4!) \int \phi^4} d\mu(\phi) \quad (\text{I.3.4})$$

To progress further towards well defined formulae, we have to introduce ultraviolet and infrared cutoffs. We will use as infrared cutoff a finite volume box  $\Lambda$ , usually a  $d$ -dimensional cube with some set of prescribed boundary conditions  $X$ . Periodic ( $X = p$ ) or Dirichlet ( $X = D$ ) boundary conditions are the most usual. The thermodynamic limit  $\Lambda \rightarrow \infty$  will be obtained through a sequence of such boxes. In this book the models considered are usually in their “high temperature phase”, in which case the thermodynamic limit, if it exists, will be independent of the particular sequence of boxes and of boundary conditions chosen. The propagator within a box, the corresponding gaussian measure and normalization factors are

noted  $C_{X,\Lambda}$ ,  $d\mu_{X,\Lambda}$ ,  $Z_{X,\Lambda}$  or simply  $C$ ,  $d\mu$ ,  $Z$ , depending on context; we try to forget subscripts or superscripts when it seems harmless.

One of the most convenient ultraviolet cutoff is the “ $\alpha$ -space cutoff”. It may also be called a heat kernel regularization of the propagator, or a regularization of the proper time of the path in the Wiener representation of the propagator as an integral over random paths. It suppresses in a smooth way the high frequencies in (I.3.2). To define it we write the  $\alpha$  or parametric representation of the propagator:

$$\hat{C}(p) = \frac{1}{(2\pi)^d} \int_0^\infty e^{-\alpha(p^2+m^2)} d\alpha \quad (\text{I.3.5})$$

$$\begin{aligned} C(x,y) &= \frac{1}{(2\pi)^d} \int_0^\infty d\alpha \int e^{ip \cdot (x-y) - \alpha(p^2+m^2)} d^d p \\ &= \frac{1}{(4\pi)^{d/2}} \int_0^\infty \frac{d\alpha}{\alpha^{d/2}} e^{-\alpha m^2 - |x-y|^2/(4\alpha)} \end{aligned} \quad (\text{I.3.6})$$

(remark that (I.3.6) is well defined except at coinciding points  $x = y$ , where for  $d \geq 2$  it is a divergent integral). We suppress the contributions of parameters  $\alpha$  less than  $\kappa$  and get:

$$C_\kappa(x,y) = \frac{1}{(4\pi)^{d/2}} \int_\kappa^\infty \frac{d\alpha}{\alpha^{d/2}} e^{-\alpha m^2 - |x-y|^2/(4\alpha)} \quad (\text{I.3.7})$$

$$\hat{C}_\kappa(p) = \frac{1}{(2\pi)^d} \int_\kappa^\infty e^{-\alpha(p^2+m^2)} d\alpha = \frac{1}{(2\pi)^d} \frac{1}{p^2 + m^2} e^{-\kappa(p^2+m^2)} \quad (\text{I.3.8})$$

When  $\kappa \rightarrow 0$ , one recovers the full propagator. In contrast with (I.3.6), (I.3.7) is well defined everywhere. It is then easy to make rigorous sense of the measure with both infrared and ultraviolet cutoffs:

$$\frac{1}{Z(X,\Lambda,\kappa)} e^{-(g/4!) \int_\Lambda \phi^4} d\mu_{X,\Lambda,\kappa}(\phi) \quad (\text{I.3.9})$$

since for  $\text{Re } g > 0$ , the exponential is bounded and in  $L^1(d\mu_{X,\Lambda,\kappa})$ . The sample fields for this cutoff measure are indeed  $C^\infty$ , and there is no difficulty to define the fourth power in (I.3.9) [Er1].

Let us also introduce the lattice regularization and the lattice  $\phi^4$  model. This is conceptually the simplest regularization scheme. Continuous space time  $\mathbb{R}^d$  is replaced by a discrete lattice grid of spacing  $\delta$ , hence by  $\delta\mathbb{Z}^d$ . We have to substitute a discrete analogue for the laplacian. We keep in the Lagrangian the term  $-\sum_{i=1}^d (\partial_i \phi)^2(x)$  but  $\partial_i$  is now the discrete derivative along the unit lattice vector  $e_i$ :

$$\partial_i \phi(x) = \frac{1}{\delta} [\phi(x + \delta e_i) - \phi(x)] \quad (\text{I.3.10}).$$

It is a straightforward computation to check that with this choice, the free field theory on the lattice is again given by a gaussian measure  $d\mu_{\text{lattice},\delta}$  with a propagator:

$$C_{\text{lattice}}^\delta = \frac{1}{(2\pi)^d} \int_{-\pi/\delta}^{\pi/\delta} d^d p \frac{e^{ip \cdot (x-y)}}{m^2 + 2\delta^{-2} \sum_{i=1}^d (1 - \cos \delta p_i)} \quad (\text{I.3.11})$$

With these notations the measure for the lattice  $\phi_d^4$  theory in a volume  $\Lambda$  is just:

$$\frac{1}{Z(\delta, \Lambda)} e^{-\frac{(g/4!)}{\delta} \sum_{x \in \Lambda \cap \delta \mathbb{Z}^d} \phi^4(x)} d\mu_{\text{lattice}, \delta}(\phi) \quad (\text{I.3.12})$$

In a finite volume, since there is only a finite number of lattice points, we can even make sense out of our starting formula (I.3.1). The lattice Laplacean built with lattice derivatives (I.3.10) involves a diagonal piece plus nearest neighbor terms, and we can therefore recast (I.3.12) in the form [GaRi]:

$$\frac{1}{Z(\delta, \Lambda)} e^{-\frac{(g/4!)}{\delta} \sum_x \phi^4(x) - \frac{(\mu^2/2)}{\delta} \sum_x \phi^2(x) + \beta \sum_{x,y} \phi(x)\phi(y)} \prod_x d\phi(x) \quad (\text{I.3.13})$$

where the sums  $\sum_x$  and the product  $\prod_x$  are taken over  $\Lambda \cap \delta \mathbb{Z}^d$  and  $\sum_{x,y}$  is performed over the pairs of nearest neighbors in  $\Lambda \cap \delta \mathbb{Z}^d$ . Furthermore:

$$\beta = \delta^{d-2} \quad \text{and} \quad \mu^2 = m^2 + 2d\delta^{-2} \quad (\text{I.3.14})$$

It is now obvious that such an expression can be taken literally. There is only a finite product of Lebesgue measures in (I.3.13), hence the functional integral is an ordinary finite dimensional integral. Every sum is finite, and for  $\operatorname{Re} g \geq 0$  this functional integral is finite, hence normalizable. A main interest of (I.3.13) is that it exhibits the lattice  $\phi^4$  theory as a finite system of classical continuous “spins”, with nearest neighbor interaction, inverse temperature  $\beta$  and single spin measure

$$e^{-(g/4!) \phi^4 - (\mu^2/2) \delta^d \phi^2} d\phi \quad (\text{I.3.15})$$

There are many techniques of statistical mechanics which apply to such a system, in particular correlation inequalities (see e.g.[Si1][Aiz][BFS][Frö]). One drawback, however, is that the form of the propagator (I.3.11) is not so simple. Hence in the next section on perturbative renormalization we will not use this lattice regularization at all.

Nevertheless when turning to constructive theory, it is good to keep several types of cutoffs at hand, and to show unicity of the limit obtained with these different cutoffs. This unicity may follow from a theorem relating the limit in a unique way to some well defined perturbation expansion, like the Borel summability theorems which will be discussed later. The advantage is that the final theory will typically retain all properties preserved by at least one set of these different cutoffs. This idea was applied to check the axioms for the Gross-Neveu model in [FMRS5]. From this point of view, we remark that the two cutoffs introduced in this section, the  $\alpha$ -space and lattice cutoff are complementary, because they violate different axioms. The  $\alpha$ -space cutoff, like the Pauli-Villars and other “momentum” space cutoffs, preserves Euclidean invariance but violates O.S. positivity. In contrast the lattice cutoff preserves a discrete version of O.S. positivity, namely reflection positivity with respect to hyperplanes of symmetry of the lattice; in the case of a cubic lattice this means the planes passing through the sites or halfway between. However the lattice is only invariant with respect to a few

discrete Euclidean symmetries, hence the lattice regularization is not Euclidean invariant.

Returning to our basic objects of study, the Schwinger functions, we observe that they still contain some rather trivial substructure which one may want to trim out in order to arrive at more interesting irreducible objects. It is a standard useful construction to build from the Schwinger functions two other classes of functions called respectively the connected Schwinger functions and the one-particle irreducible (in short 1PI) Schwinger functions (in statistical mechanics connected functions are called Ursell functions or cumulants).

The connected Schwinger functions are given by:

$$C_N(z_1, \dots, z_N) = \sum_{P_1 \cup \dots \cup P_k = [1, N]; P_i \cap P_j = \emptyset} (-1)^{k+1} \prod_{i=1}^k S_{p_i}(z_{j_1}, \dots, z_{j_{p_i}}) \quad (\text{I.3.16})$$

where the sum is performed over all distinct partitions of  $[1, N]$  into  $k$  subsets  $P_1, \dots, P_k$ ,  $P_i$  being made of  $p_i$  elements called  $j_1, \dots, j_{p_i}$ . For instance the connected 4-point function in the high temperature region of the  $\phi^4$  theory, where all odd Schwinger functions vanish due to the unbroken  $\phi \rightarrow -\phi$  symmetry, is simply given by:

$$\begin{aligned} C_4(z_1, \dots, z_4) &= S_4(z_1, \dots, z_4) - S_2(z_1, z_2)S_2(z_3, z_4) \\ &\quad - S_2(z_1, z_3)S_2(z_2, z_4) - S_2(z_1, z_4)S_2(z_2, z_3) \end{aligned} \quad (\text{I.3.17})$$

The 1PI functions, also called vertex functions,  $\Gamma_N(z_1, \dots, z_N)$  are slightly harder to derive from the Schwinger functions. Nevertheless 1PI functions may be defined rigorously directly from the Schwinger functions, without relying on perturbation theory as their names seem to suggest. This direct construction involves combining the formalism of the first Legendre transform with T. Spencer's idea of testing irreducibility through decoupling surfaces [Sp2]. For a complete presentation of this point of view we refer to [CFR]. Here we will give simply a naive formal definition, and admit that vertex functions can be derived rigorously from Schwinger functions and vice versa. We will also admit that the perturbation theory of these objects coincide with the usual graphical definition of one particle irreducibility.

The generating functional for the connected functions is:

$$C(J) = \text{Log} Z(J) = \text{Log} \int e^{\phi \cdot J} d\nu(\phi) \quad (\text{I.3.18})$$

in the sense that

$$C_N(z_1, \dots, z_N) = \prod_{i=1}^N \frac{\delta}{\delta J(z_i)} C(J)|_{J=0} \quad (\text{I.3.19})$$

The first Legendre transform is defined by inverting

$$A(J)(x) = \frac{\delta}{\delta J(x)} C(J) - \frac{\delta}{\delta J(x)} C(J)|_{J=0} \quad (\text{I.3.20})$$

$$\Gamma(A) = C\{J(A)\} - \frac{\delta C}{\delta J}\{J(A)\} \cdot J(A) \quad (\text{I.3.21})$$

where the “scalar product” is defined by  $\frac{\delta C}{\delta J} \cdot J = \int dx \frac{\delta C}{\delta J(x)} J(x)$

This Legendre transform is the generating functional for the vertex functions, which are therefore given by:

$$\Gamma_N(z_1, \dots, z_N) = \prod_{i=1}^N \frac{\delta}{\delta A(z_i)} \Gamma(A)|_{A=0} \quad (\text{I.3.22})$$

## I.4 Feynman graphs and amplitudes

### A) Graphs

The perturbation expansion is an expansion in powers of the coupling constant. In the case of the Schwinger functions (I.3.4) this means that one writes:

$$S_N(z_1, \dots, z_N) = \frac{1}{Z} \sum_{n=0}^{\infty} \frac{(-g)^n}{n!} \left[ \int \frac{\phi^4(x)}{4!} \right]^n \phi(z_1) \dots \phi(z_N) d\mu(\phi) \quad (\text{I.4.1})$$

It is possible to perform explicitly the functional integral of a polynomial in the fields with respect to a gaussian measure. The result, called in physics “Wick’s theorem”, gives at any order a sum over “Wick contractions”, i.e. ways of pairing together the fields in (I.4.1). More precisely, we can label the integration variables in (I.4.1) as  $x_1, \dots, x_n$ , and the 4 fields in a monomial  $\phi^4$  as  $\phi_1 \phi_2 \phi_3 \phi_4$ . Then we commute the spatial integral over  $x_1, \dots, x_n$  and the functional integral. The result is a gaussian integral over a polynomial in  $\phi$  of degree  $4n + N$ , namely  $\prod_{i=1}^n \prod_{k=1}^4 \phi^k(x_i) \prod_{j=1}^N \phi(z_j)$ . By the rules of gaussian integration, the result is:

$$\sum_{\Pi} \prod_{l \in \Pi} C(x_l, y_l) \quad (\text{I.4.2})$$

where the sum is over contraction schemes  $\Pi$  namely the partitions of the set of all  $4n + N$  fields into  $2n + N/2$  pairs  $l$ , called lines. For such a line  $l$ ,  $x_l$  and  $y_l$  are the arguments of the fields in the pair. (There is no ordering ambiguity in (I.3.2) since  $C(x, y) = C(y, x)$ ). There are exactly  $(4n + N - 1)(4n + N - 3)\dots 5.3.1 = (4n + N - 1)!!$  such contraction schemes.

Formally at order  $n$  the result of perturbation theory is therefore simply the sum over all these schemes  $\Pi$  of the spatial integrals over  $x_1, \dots, x_n$  of the integrand (I.3.2) times the factor  $\frac{1}{n!} \left( \frac{-g}{4!} \right)^n$ . These integrals are then functions (in fact distributions) of the external positions  $z_1, \dots, z_N$ . But they may diverge either because they are integrals over all of  $\mathbb{R}^4$  (no volume cutoff) or because of the singularities in the propagator  $C$  at coinciding points.

For practical computations, it is obviously more convenient to gather all the contractions which lead to the same topological structure, hence the same integral. This leads to the notion of Feynman graphs or diagrams. To any such graph is associated a contribution or amplitude, which is the sum of the contributions associated with the corresponding set of Wick contractions. The Feynman rules summarize how to compute this amplitude with its correct combinatoric factor. In the case of Euclidean  $\phi_d^4$  these rules are rather simple and we will describe them in detail, introducing some terminology and notations which are used throughout the rest of the book.

A Feynman diagram is a set of labelled vertices together with a set of lines between them. Each line is hooked at each end to a given vertex. Such an “end” of a line is called a half-line. Half lines are the remnants of the fields after Wick’s theorem has been applied. There may be vertices called the external vertices, labelled as  $V_1, \dots, V_N$ , to which only a single half-line is hooked. And there may be internal vertices, labelled  $v_1, \dots, v_n$ , to which more than one half-line is hooked. We call the diagram a  $\phi^4$  diagram if there are exactly four half-lines hooked to any internal vertex. Examples of  $\phi^4$  diagrams are shown in Fig.I.4.1. Clearly

the (bare)  $\phi^4$  perturbation theory leads to  $\phi^4$  diagrams whose  $N$  external vertices correspond to the external fields in (I.3.1). For the renormalization of  $\phi_4^4$  it will be convenient to define also generalized  $\phi^4$  diagrams as diagrams with four types of internal vertices, the regular internal vertices with four half lines plus three kinds of special internal vertices for counterterms:

- the coupling constant counterterm, with four half lines,
- the mass counterterms, with two half-lines,
- the wave function counterterms, again with two half-lines, and two arrows to distinguish them from the mass counterterms. We will see that these arrows correspond to derivative acting on the corresponding propagators.

The counterterms are pictured as thicker “blobs” to distinguish them from the ordinary ones. Generalized  $\phi^4$  diagrams are shown in Fig.I.4.2. Usually it will be clear from context which class of graphs is considered.

Feynman diagrams without any labelling of the internal vertices are called Feynman graphs. Therefore there are at most  $n!$  diagrams corresponding to a graph.

We use always the following notations for a graph  $G$ :

- $n(G)$  or simply  $n$  is the number of internal vertices of  $G$ , or the order of the graph.
- $l(G)$  or  $l$  is the number of internal lines of  $G$ , i.e. lines hooked at both ends to an internal vertex of  $G$ . We use also widely  $l$  as an index for lines; in cases where confusion is possible, we try to note always the number of lines as  $l(G)$ .
- $N(G)$  or  $N$  is the number of external vertices of  $G$ ; it corresponds to the order of the Schwinger function one is looking at. When  $N = 0$  the graph is a vacuum graph, otherwise it is called an  $N$ -point graph.
- $c(G)$  or  $c$  is the number of connected components of  $G$ ,
- $L(G)$  or  $L$  is the number of independent loops of  $G$ .

For a *regular*  $\phi^4$  graph, i.e. a graph which has no line hooked at both ends to external vertices, we have the relations:

$$l(G) = 2n(G) - N(G)/2 \quad (\text{I.4.3})$$

$$L(G) = l(G) - n(G) + c(G) \quad (\text{I.4.4})$$

and, for a connected graph:

$$L(G) = n(G) + 1 - N(G)/2 \quad (\text{I.4.5})$$

We like to define the superficial degree of convergence, which is the opposite of the more standard degree of divergence. For  $\phi_d^4$  it is:

$$\omega(G) = 2l(G) - dL(G) \quad (\text{I.4.6})$$

so that for a connected graph:

$$\omega(G) = (4 - d)n(G) + \frac{d - 2}{2}N(G) - d \quad (\text{I.4.7})$$

It will be important also to define what we call a subgraph; this is not a completely straightforward notion. A *subgraph*  $F$  of a graph  $G$  is a subset of

internal lines of  $G$ . Hence there are exactly  $2^{l(G)}$  subgraphs in  $G$ . We call the lines in the subset defining  $F$  the internal lines of  $F$ , and their number is simply  $l(F)$ , as before. Similarly all the vertices of  $G$  hooked to at least one of these internal lines of  $F$  are called the internal vertices of  $F$  and considered to be in  $F$ ; their number by definition is  $n(F)$ . But remark that no external vertex of  $G$  can be of this kind. Precisely for this reason, the notion of external vertices does not generalize simply to subgraphs. Nevertheless for power counting we need at least to define a generalization of the number  $N$  for subgraphs. A good convention is to call external half-line of  $F$  every half-line of  $G$  which is not in  $F$  but which is hooked to a vertex of  $F$ ; it is then the number of such external half-lines which we call  $N(F)$ . With this convention one has for  $\phi^4$  subgraphs the same relation (I.4.3) as for regular  $\phi^4$  graphs.

The definitions of  $c$ ,  $L$  and  $\omega$  for subgraphs are then straightforward, and relations (I.4.4)-(I.4.7) extend to them.

Finally it will be useful to have a particular name for the vertices of  $F$  attached to at least one external half-line of  $F$ . We call them the “border” vertices. The internal vertices of  $F$  which are not border-vertices are called inside-vertices. We will see that border-vertices partly play the role of external vertices for subgraphs, in the sense that it is sometimes natural to integrate over the position of inside-vertices and to consider the corresponding “subamplitudes” as functions (in fact distributions) of the positions of the border-vertices. Examples of subgraphs are pictured in Fig. I.4.3.

## B) Amplitudes

To compute the amplitude associated to a  $\phi^4$  diagram, we have to add the contributions of the corresponding contraction schemes. This is summarized by the rules:

- To each line  $l_j$  with end vertices at positions  $x_j$  and  $y_j$ , associate a propagator  $C(x_j, y_j)$ .
- To each internal vertex, associate  $(-g)/4!$ .
- Count all the contraction schemes giving this diagram. The number should be of the form  $(4!)^n/S_1(G)$  where  $S_1(G)$  is an integer called the symmetry factor of the diagram. (If it is not the case, try again). The  $4!$  represents the permutation of the fields hooked to an internal vertex. Remark that  $S_1(G)$  does not depend on the labelling of internal vertices, hence is a function of the underlying graph.
- Multiply all these factors, divide by  $n!$  and sum over the position of all internal vertices.

This gives the bare amplitude of a diagram, neglecting possible divergences. To get the bare amplitude for a *graph*, we count the number of diagrams giving rise to this graph (by erasing internal labelling). This number is  $n!/S_2(G)$ , where  $S_2(G)$  is an integer.  $S(G) = S_1(G)S_2(G)$  is called the total symmetry factor for the graph  $G$ . Some of these factors are shown in Fig. I.4.4.

The formula for the bare amplitude of a graph is therefore, as a distribution in  $z_1, \dots, z_N$ :

$$A_G(z_1, \dots, z_N) \equiv \int \prod_{i=1}^n dx_i \prod_{l \in G} C(x_l, y_l) \quad (\text{I.4.8})$$

As stated above, this integral suffers of possible divergences. But the corresponding quantities with both volume cutoff and ultraviolet cutoff  $\kappa$  are well defined. They are:

$$A_{G,\Lambda}^\kappa(z_1, \dots, z_N) \equiv \int_{\Lambda^n} \prod_{i=1}^n dx_i \prod_{l \in G} C_\kappa(x_l, y_l) \quad (\text{I.4.9})$$

The integrand is indeed bounded and the integration domain is a compact box.

Returning to (I.4.1) the *unnormalized* Schwinger functions are formally given by the sum over all graphs with the right number of external lines of the corresponding Feynman amplitudes:

$$ZS_N = \sum_{\phi^4 \text{ graphs } G \text{ with } N(G)=N} \frac{(-g)^{n(G)}}{S(G)} A_G \equiv \sum_n (-g)^n b_n^N \quad (\text{I.4.10})$$

$Z$  itself, the normalization, is given by the sum of all vacuum amplitudes:

$$ZS_N = \sum_{\phi^4 \text{ graphs } G \text{ with } N(G)=0} \frac{(-g)^{n(G)}}{S(G)} A_G \equiv \sum_n (-g)^n b_n^0 \quad (\text{I.4.11})$$

From translation invariance, we do not expect  $A_{G,\Lambda}^\kappa$  to have a limit as  $\Lambda \rightarrow \infty$  if there are vacuum subgraphs in  $G$ . But we can remark that an amplitude factorizes as the product of the amplitudes of its connected components at least as far as the integral is concerned; the only subtle point is the combinatoric factors, which are not easy to disentangle at the level of graphs because of the factor  $S(G)$ , but which are reasonably transparent at the level of contraction schemes. Let us consider (I.4.11). Clearly each contraction scheme  $w$  creates a partition of the  $n$  vertices of (I.4.1) into  $k$  subsets  $V_1, \dots, V_k$  of  $n_1, \dots, n_k$  vertices which are the connected components of the contraction scheme. Summing first over the partition and then over the schemes which give rise to this partition, one has to take care of the proper combinatoric factor  $n! / \prod_{i=1, \dots, k} n_i!$  to attribute the vertices to each

subset, and of an additional  $\frac{1}{k!}$  since partitions are unordered sums; then the initial sum is put in the form  $\frac{1}{k!} \sum_{V_1, \dots, V_k} \prod_i A(V_i)$  where  $A(V)$  is a power series exactly similar to the initial one, except that when applying Wick's theorem, only Wick's contractions which connect together all the vertices of  $V$  are allowed. Therefore one has proved that the pressure  $p(\Lambda) = \frac{1}{|\Lambda|} \log Z(\Lambda)$  has a perturbation expansion given by the sum of all vacuum *connected* Feynman amplitudes:

$$p(\Lambda) = \sum_{\phi^4 \text{ connected vacuum graphs } G} \frac{1}{|\Lambda|} \frac{(-g)^{n(G)}}{S(G)} A_{G,\Lambda} \equiv \sum_n (-g)^n a_n^0 \quad (\text{I.4.12})$$

Later in this book we introduce the Mayer expansion which is a systematic way of computing the logarithm of partition functions, and we invite the reader to prove again equation (I.4.12) using this general method.

When  $G$  is a vacuum connected graph, there is a single overall translation invariance in its amplitude; therefore using the rapid decay of (I.3.6) for large  $|x - y|$  we expect  $\frac{1}{|\Lambda|} A_{G,\Lambda}$  to have a limit as  $\Lambda \rightarrow \infty$ , hence the pressure, not the normalization, is the correct quantity to consider in the thermodynamic limit. For  $G$  such a vacuum connected graph we write:

$$\lim_{\Lambda \rightarrow \infty} \frac{1}{|\Lambda|} A_{G,\Lambda} = I_G \equiv \int_{\mathbb{R}^d} \dots \int_{\mathbb{R}^d} \prod_{i \neq 1} dx_i \prod_l C(x_l, y_l) |_{x_1=0} \quad (\text{I.4.13})$$

With similar simple combinatoric arguments or with the formalism of the Mayer expansion, we can also factorize the vacuum graphs in the expansion of unnormalized Schwinger functions, and get for the normalized ones:

$$S_N = \sum_{\substack{\phi^4 \text{ graphs } G \text{ with } N(G)=N \\ G \text{ without any vacuum subgraph}}} \frac{(-g)^{n(G)}}{S(G)} A_G \equiv \sum_n (-g)^n a_n^N \quad (\text{I.4.14})$$

Again in (I.4.14) it is possible to pass to the thermodynamic limit (in the sense of formal power series) because using the exponential decrease of the propagator, each individual graph has a limit (at fixed external arguments). There is of course no need to divide by the volume for that because each connected component in (I.4.14) is tied to at least one external source, and they provide the necessary breaking of translation invariance.

Finally with the help of the Mayer expansion or by direct reasoning one can determine the perturbative expansions for the connected Schwinger functions and the vertex functions. As expected the connected Schwinger functions are given by sums over connected amplitudes:

$$C_N = \sum_{\substack{\phi^4 \text{ connected graphs } G \text{ with } N(G)=N}} \frac{(-g)^{n(G)}}{S(G)} A_G \equiv \sum_n (-g)^n a_{n,c}^0 \quad (\text{I.4.15})$$

and the vertex functions are the sums of the *amputated* amplitudes for 1PI or proper graphs. They are the graphs which remain connected even after removal of any given internal line. The amputated amplitudes are defined in momentum space by omitting the Fourier transform of the propagators of the external lines. It is therefore convenient to write these amplitudes in the so called momentum representation:

$$\Gamma_N = \sum_{\substack{\phi^4 \text{ proper graphs } G \text{ with } N(G)=N}} \frac{(-g)^{n(G)}}{S(G)} A_G^T(z_1, \dots, z_N) \equiv \sum_n (-g)^n a_{n,p}^N \quad (\text{I.4.16})$$

$$A_G^T(z_1, \dots, z_N) \equiv \frac{1}{(2\pi)^{dN/2}} \int dp_1 \dots dp_N e^{i \sum p_i z_i} A_G(p_1, \dots, p_N) \quad (\text{I.4.17})$$

$$A_G(p_1, \dots, p_N) = \int \prod_{l \in G} \frac{d^d p_l}{p_l^2 + m^2} \prod_{v \in G} \delta\left(\sum_l \epsilon_{v,l} p_l\right) \quad (\text{I.4.18})$$

Remark in (I.4.18) the  $\delta$  functions which ensure momentum conservation at each internal vertex  $v$ ; the sum inside is over both internal and external lines; each internal line is oriented in an arbitrary way and each external line is oriented towards the inside of the graph; then the incidence matrix  $\epsilon(v, l)$  is 1 if the line  $l$  arrives at  $v$ , -1 if it starts from  $v$  and 0 otherwise. Remark also that there is an overall momentum conservation rule  $\delta(p_1 + \dots + p_N)$  hidden in (I.4.18). By Euclidean invariance it is clear that  $I_G$  depends only of the Euclidean invariants built on the external momenta, which are scalars; this remark is important for interpolation of amplitudes to non-integer dimension [Er2]. Notice that in the  $\phi^4$  theory a connected vacuum graph is always proper, and indeed for  $N = 0$  the definition (I.4.18) coincides with (I.4.13). We may define  $P_v$  as the sum of all external momenta entering vertex  $v$ , in which case the momentum conservation function at each vertex becomes  $\delta(P_v + \sum_l \epsilon_{v,l} p_l)$ , where the sum runs now over internal lines  $l$  only.

The drawback of the momentum representation lies in the necessity for practical applications to eliminate the  $\delta$  functions by a “momentum routing” prescription, and there is no canonical choice for that. An other interesting representation is the parametric representation which is obtained after the  $x$  space or  $p$  space integration has been explicitly performed, using respectively representations (I.3.6) or (I.3.5) of the propagator. This is possible because these integrations are quadratic. The result is a compact formula with one scalar integration over a parameter  $\alpha$  for each internal line of the graph. To write down explicitly this formula requires (for the first time in this book) the key combinatoric notion of a tree, so we pause for a while to give some corresponding definitions and notations.

### C) Trees

A tree can be defined as a graph which is connected and without loops. If the constraint of connectedness is removed we have a forest, which is therefore nothing more than a finite set of trees. Like in the distinction between graphs and diagrams, we have to decide whether the vertices of the tree will be labeled or not. One of the most useful aspects of trees lies in the partial ordering relations they provide. More precisely, to any particular vertex of a tree is associated a natural partial ordering obtained by drawing the tree with that particular vertex as its root. This partial ordering means that a vertex is “higher” than an other if the particular unique path connecting it to the root passes through the other. Maximal vertices may be called the “leaves” and the tree together with its particular vertex a “rooted tree”.

Up to now this description is rather intuitive; sometimes the problem of counting trees arise and the reader may like to have a more set-theoretic definition. We propose to define both unordered and ordered trees. A labeled unordered tree (or simply a tree) with  $n$  vertices is simply a set  $T$  of  $n-1$  unordered pairs or links  $\{l_1, \dots, l_{n-1}\}$  among  $n$  elements (vertices) which can be represented as  $\{1, \dots, n\}$ , with the property (P):

For any couple of vertices  $(i, j)$  there is a *unique* path from  $i$  to  $j$  in  $T \quad (\text{P})$

A rooted tree is a tree plus the choice of a particular vertex; hence there are  $n$  rooted trees with  $n$  vertices for any single tree.

One may call “tree shapes” the quotient of trees by the equivalence relation of erasing the labels of vertices. For  $n = 3, 4, 5$ , there are respectively 1, 2 and 3 tree shapes, but unfortunately such a simple series does not continue and we do not know whether there is a simple formula for the number of tree shapes at any order.

The tree shape is somewhat probed by the coordination numbers of the tree, which are the numbers  $d_i$  of pairs containing  $i$ . It is for the number of labeled unordered trees, with or without fixed coordination numbers, that Cayley’s theorem gives a simple formula:

#### **Theorem I.4.1: Cayley’s theorem**

The number of labeled unordered trees with  $n$  vertices is  $n^{n-2}$ . The number of such trees with fixed coordination numbers  $d_i$  is  $\prod_i \frac{(n-2)!}{(d_i-1)!}$ .

For instance for  $n = 4$ , there are sixteen trees, depicted in Fig.I.4.5. For coordination numbers  $(3,1,1,1)$  there is a single possibility, because the tree shape is the first one and the vertex with  $n_i = 3$  must be at the center; for coordination numbers  $(2,2,1,1)$ , there are 2 possibilities.

It is also sometimes useful to define a “labeled ordered rooted tree” or in short an ordered tree as a rooted tree plus a total ordering which is compatible with the partial ordering relation of the rooted tree, and such that the root has label 1. Hence it may be viewed as a permutation  $j$  of  $\{1, \dots, n\}$  satisfying  $j(1) = 1$ , plus an “ancestor” function  $a$  from  $\{1, \dots, n-1\}$  to  $\{1, \dots, n-1\}$  which satisfies, instead of (P), the property:

$$A(k) \leq k \quad \forall k \quad (\text{P}')$$

The links of the trees are the pairs  $(i(k), j(k+1))$ ,  $k = 1, \dots, n-1$ , where  $i(k) \equiv j(a(k))$ . Such a structure appears naturally in the cluster expansion scheme of section III.1. There are at most  $(n-1)!$  ordered trees for an unordered one; there are for instance 176 ordered trees for  $n = 4$ , but no simple general formula.

Inside a graph one defines a spanning tree (or in short, a tree) as a subset of internal lines which is a tree and which connects together all the internal vertices of the graph (hence is maximal for the number of lines); a two tree is a spanning tree minus one line, hence it splits the internal vertices into two connected subsets. Remark that graphically a two tree is either a forest of two trees or a single tree plus an isolated vertex (if the line removed from the spanning tree was an extremal one).

#### **D) Parametric representation**

With these notions, let us return to the parametric representation and prove:

$$A_G(p_1, \dots, p_N) = \delta(\sum_v P_v) \int_0^\infty \prod_l d\alpha_l e^{-\sum_l \alpha_l m^2 - V_G(\alpha, p)/U_G(\alpha)} \frac{1}{[U_G(\alpha)]^{d/2}} \quad (\text{I.4.19})$$

where  $U_G$  and  $V_G$  are polynomials in  $\alpha$  depending on the particular topology of the graph, called the Symanzik polynomials. Their explicit expression is:

$$U_G = \sum_S \prod_{l \text{ not in } T} \alpha_l \quad (\text{I.4.20})$$

$$V_G(p, \alpha) = \left( \sum_T \prod_{l \text{ not in } T} \alpha_l \right) \left( \sum_{a \in T_1} p_a \right)^2 \quad (\text{I.4.21})$$

In (I.4.20) the sum runs over the spanning trees  $S$  of  $G$ , and in (I.4.21) over the two-trees  $T$  of  $G$  which separate  $G$  into two connected components, each containing a non empty set of external lines, one of which is  $T_1$  (by overall momentum conservation, (I.4.21) does not change if  $T_1$  is replaced by the set of external lines of the other connected component, which is the complementary of  $T_1$ ).

**Proof** Following [IZ], we rewrite the  $\delta$  functions expressing momentum conservation as:

$$(2\pi)^d \delta(P_v + \sum_l \epsilon_{v,l} p_l) = \int dy_v e^{-iy_v \cdot (p_v + \sum_l \epsilon_{v,l} p_l)} \quad (\text{I.4.22})$$

Then we exchange the order of integration and integrate over each  $p_l$  (recall that this is just a formal computation of bare amplitudes):

$$\int d^d p_l e^{-\alpha_l p_l^2 - i \sum_v y_v \epsilon_{v,l} p_l} = \frac{1}{(\prod_l \alpha_l)^{d/2}} e^{-\left(\sum_v y_v \epsilon_{v,l}\right)^2 / 4\alpha_l} \quad (\text{I.4.23})$$

Therefore:

$$A_G(p_1, \dots, p_N) = \int \prod_v \frac{dy_v}{(2\pi)^d} e^{-iy_v \cdot P_v} \prod_l \frac{d\alpha_l}{(\pi\alpha_l)^{d/2}} e^{-\alpha_l m^2 - \left(\sum_v y_v \epsilon_{v,l}\right)^2 / 4\alpha_l} \quad (\text{I.4.24})$$

To integrate over  $y$ 's variables, we shift by  $y_n$ , the last variable, defining  $z_v = y_v + y_n$  for  $v < n$ , and  $z_n = y_n$ . The jacobian is one. Since  $\sum_v \epsilon_{v,l} = 0$ , the integration over  $z_n$  simply yields the overall momentum conservation  $\delta(\sum_v P_v)$  and we are left with  $n - 1$  quadratic integrations. We define the  $n - 1$  by  $n - 1$  square matrix  $[d_G]_{i,j}$  for  $i \leq n - 1$ ,  $j \leq n - 1$  as

$$d_G(\alpha)_{i,j} = \sum_l \epsilon_{i,l} \frac{1}{\alpha_l} \epsilon_{j,l} \quad (\text{I.4.25})$$

The gaussian integration in (I.4.24) corresponds to

$$e^{-(\frac{1}{4})(y+2iPd_G^{-1})d_G(y+2iPd_G^{-1})^t - Pd_G^{-1}P} \quad (\text{I.4.26})$$

with obvious matrix notations, and the result is  $\det(d_G^{-d/2} e^{-Pd_G^{-1}P})$ . Combining this with the factor  $\prod_l \alpha_l^{-d/2}$  in (I.4.24) achieves the proof of (I.4.19)-(I.4.21), provided one can show that  $d_G$  is non singular and that

$$\det d_G(\alpha) = \sum_T \prod_{l \in T} \frac{1}{\alpha_l} = U_G(\alpha) \quad (\text{I.4.27})$$

$$\sum_{i,j=1}^{n-1} P_i [d_G^{-1}]_{ij} P_j = \frac{V_G(p, \alpha)}{U_G(\alpha)} \quad (\text{I.4.28})$$

These results are called topological formulas. A detailed study of such formulas is in [Nak]. We will limit ourselves to a quick proof of (I.4.27), using the Binet-Cauchy formula. This formula computes the determinant of a  $p$  by  $p$  matrix  $C$  which is the product of two rectangular matrices  $A$  and  $B$ , the first  $p$  by  $q$  and the second  $q$  by  $p$ . We may assume  $p < q$  for non triviality and the formula is simply:

$$\det C = \sum_M \det A_M \det B_M \quad (\text{I.4.29})$$

where  $M$  runs over all subsets of  $p$  indices among  $q$ , and the minors  $A_M$  and  $B_M$  are obtained by deleting in  $A$  (respectively in  $B$ ) the columns (respectively the lines) with indices not in  $M$ . Now let us call  $\epsilon_n$  the reduced incidence matrix obtained from  $\epsilon$  by deleting the  $n$ -th line (corresponding to vertex  $n$ ). We have  $d_G = \epsilon_n \alpha \epsilon_n^t$ , where  $\alpha$  is the diagonal matrix  $\alpha_l \delta_{ll}$ . From (I.4.29), (I.4.27) obviously follows if we can check that  $\epsilon_n$  has rank  $n - 1$  and that its  $n - 1$  by  $n - 1$  minors  $\epsilon_{n,M}$  are  $\pm 1$  or 0 depending on whether the lines kept in  $M$  form a tree or not. If in  $\epsilon_n$  we keep a subset of lines  $M$  which is not a tree, it has to contain a closed circuit, and the sum of incidence numbers along such a circuit is 0; this gives a linear relation which proves that the corresponding minor is 0. Finally when  $M$  is a tree, the fact that the minor is  $\pm 1$  can be checked by induction on the number of vertices of  $G$ ; we consider a line  $l$  of  $M$  hooked to the vertex  $n$  deleted in  $\epsilon_n$ ; then the corresponding column in  $\epsilon_n$  has only one nonzero element  $[\epsilon_n]_{il}$ . Expanding the corresponding minor with respect to that column we obtain up to a sign the same problem for a smaller graph  $G/l$  in which line  $l$  has been deleted and vertices  $n$  and  $i$  are collapsed. Indeed  $M - l$  is still a spanning tree of  $G/l$  and we may assume the collapsed vertex to be the distinguished one in  $G/l$  so that the induction hypothesis applies.

The parametric representation is perhaps the most elegant of all and is well suited for the study of particular amplitudes, and also for general results like analyticity properties, existence and nature of asymptotic expansions in various regimes, and also for writing renormalization operators which act directly on the integrand [BL][BZ]. It is in this representation that large order bounds for the renormalized  $\phi^4$  theory were first obtained [dCR1]. Nevertheless we will not use it in these notes, because constructive theory relies heavily on  $x$  space which is lost in the parametric representation.

In the next chapter we will use a mixed representation based on (I.3.6) but with a further discrete slicing of the  $\alpha$  integration. This representation which we call the multiscale representation, was introduced and used for perturbative and constructive studies in the series of papers [FMRS1-5] which is at the origin of this book. Many variations over this theme are worth being studied. For instance with the pure  $x$  and  $\alpha$  propagator (without the discrete slicing), perhaps the most elegant and dense study of the perturbative renormalization of  $\phi^4$  has been given recently in [Hu]. Nevertheless throughout this book we stick as much as possible to the original representation of [FMRS1], because at least up to now a discrete slicing seems necessary for constructive purposes.

## I.5 Borel summability

Borel summability is one of many possible substitutes for ordinary summability. It has proved very useful for the analysis of many divergent series met in physics, and provides in particular a natural framework for the study of the perturbative series met in this book. Therefore we include this brief section to recall what it means.

An analytic function inside its domain of analyticity is the sum of its Taylor series, so that all information about this function is embedded in the list of Taylor coefficients at an interior point. Ordinary summation then provides a one to one correspondence (at least inside a convergence disk) between convergent power series and analytic functions. Borel summability is a method to extend this one to one correspondence to the case of an analytic function expanded in Taylor series at a point on the *border* of the domain of analyticity. Such a Taylor series is no longer summable in the ordinary sense, but under precise conditions it may still be associated in a unique way to a function with sufficiently large analyticity domain and sufficiently strong asymptoticity to this Taylor series. In this sense, again all the information about the function (the “Borel sum”) is still embedded in the series (the “Borel series”).

In [Ha] there is a version of Borel summability which was popular among physicists, until A. Sokal found a somewhat simpler and more natural version [Sok1]. This version was in fact a rediscovery of a theorem of Nevanlinna [Ne], published in 1919:

### Theorem I.5.1 (Nevanlinna-Sokal)

Let  $f$  be analytic in the disk  $C_R = \{y | \operatorname{Re} y^{-1} > 1/R\}$  of Fig.I.5.1. Suppose  $f$  admits an asymptotic power series  $\sum a_k y^k$  (its Taylor series at the origin) hence:

$$f(y) = \sum_{k=0}^{r-1} a_k y^k + R_r(y) \quad (\text{I.5.1})$$

such that the bound

$$|R_r(y)| \leq c\sigma^r r! |y|^r \quad (\text{I.5.2})$$

holds uniformly in  $r$  and  $y \in C_R$ , for some constants  $\sigma$  and  $C$ . Then  $f$  is Borel summable, which means that the power series  $\sum_k a_k \frac{t^k}{k!}$  converges for  $|t| < \frac{1}{\sigma}$ , that it defines a function  $B(t)$  which has an analytic continuation in the strip  $S_\sigma = \{t | \operatorname{dist}(t, \mathbb{R}^+) < \frac{1}{\sigma}\}$  of Fig I.5.1, and that this function satisfies the bound

$$|B(t)| \leq \text{const.} e^{\frac{t}{\sigma}} \quad \text{for } t \in \mathbb{R}^+ \quad (\text{I.5.3})$$

Finally  $f$  is represented by the following absolutely convergent integral:

$$f(y) = \frac{1}{y} \int_0^\infty e^{-\frac{t}{y}} B(t) dt \quad \text{for } y \in C_R \quad (\text{I.5.4})$$

Under these conditions,  $f$  is said to be Borel summable, and  $B$  is called its Borel transform. The complex  $t$  plane is called the Borel plane. There is a reciprocal to this theorem which states that starting with a given power series  $\sum a_k y^k$  if the power series  $\sum a_k \frac{t^k}{k!}$  converges in a disk  $|t| < \frac{1}{\sigma}$ , admits an analytic

continuation  $B(t)$  in the strip  $S_\sigma$  and satisfies the bound (I.5.3) in this strip, then the function  $f$  defined by the integral representation (I.5.4) is analytic in  $C_R$ , has  $\sum a_k y^k$  as Taylor series at the origin and satisfies the uniform remainder estimates (I.5.1-2). In this case we say that the series  $\sum a_k y^k$  is Borel summable, that the series  $\sum a_k \frac{t^k}{k!}$  is its Borel transform, and that the function  $f$  is its Borel sum. In conclusion Borel summable series and Borel summable functions are in correspondence just as are ordinary series and germs of analytic functions. However the analytic continuation in the Borel strip involved in the construction of the function from its series is usually untractable, so that it is the direct theorem which is used in practice.

This theorem is by no means optimal, in the sense that for many typical power series it does not reconstruct the maximal analyticity domain of the Borel sum; this can be checked even in the simplest non-trivial case  $a_k = (-1)^k k!$ . One can use modified versions to recover in this kind of situations a larger domain of analyticity; also for power series with different large order behavior, such that  $(n!)^\alpha$  for instance, the Nevanlinna-Sokal theorem should be modified (for instance by applying some conformal maps of the disk of Fig.I.6.1 onto other domains). We do not attempt a review on the extensive mathematical litterature on this subject. There are also many results of a general nature on Borel summable functions, which state for instance that products, derivatives, inverse, composition etc... of Borel summable functions remain Borel summable basically provided these operations make sense (see e.g. [AM][BCS]). Regular Borel summability is adequate for the problems discussed in this book; in particular the Borel plane is the natural setting to study the large order behavior of  $\phi^4$  (section II.6) and the relationship between asymptotically free models and their perturbation expansion (sections II.5, III.3-4).

**PART II: PERTURBATIVE RENORMALIZATION**

## II.1 The multiscale representation and a bound on convergent graphs

*Suppose we had a strictly finite theory,  
with bounded propagators, bounded integrals and all that.*

*Individual diagrams in such a theory are then bounded  
by a pure power law as a function of their order  $n$ .  
– G. 't Hooft, “Can we make sense out of QCD?”*

In this section we prove in detail a bound which we think is the best pedagogical introduction to multiscale expansions. It is a kind of “uniform” Weinberg theorem. The original Weinberg theorem [We] proves that any graph  $G$  for which each connected subgraph  $S$  has a positive superficial degree of convergence  $\omega(S)$  has indeed a finite amplitude (hence does not require renormalization). From (I.4.7) this means in the case of  $\phi_4^4$  that any graph  $G$  such that each connected subgraph  $S$  has at least 6 external lines (5 being forbidden) is finite. Such graphs will be called completely convergent; some examples of them are shown in Fig.II.1.1.

The best proof of the Weinberg theorem is in the parametric representation, by use of Hepp's sectors [He1-2]. Although the result agrees with intuition based on “power counting”, it is not trivial and rests on the particular structure of the Feynman integrand, more precisely on the structure of Symanzik's polynomials. It is of course not true that any multidimensional integral converges when each subset of parameters has convergent scaling properties; for instance the integral

$$\int_0^\infty \int_0^\infty \int_0^\infty \frac{d\alpha_1 d\alpha_2 d\alpha_3 \alpha_2 \alpha_3 e^{-\alpha_1 - \alpha_2 - \alpha_3}}{(\alpha_1^2 + \alpha_2 \alpha_3)^2} \quad (\text{II.1.1})$$

is not absolutely convergent, although it is superficially convergent under scaling of any given subset of  $\alpha_1, \alpha_2, \alpha_3$  (one could find more complicated examples, superficially more resembling to actual Feynman amplitudes). But the Symanzik polynomial  $U_G$  of the last section cannot be of this type. The key property of  $U_G$ , not shared by the example above, was called the *fine* (Factorized IN Each sector) property in [dCM]. For  $\sigma$  any permutation of  $[1, \dots, l]$ , the associated Hepp's sector  $H_\sigma$  is the region  $\alpha_{\sigma(1)} \leq \dots \leq \alpha_{\sigma(l)}$ . By definition, a fine polynomial  $P$  of  $l$  variables  $\alpha_1, \dots, \alpha_l$  is such that in each Hepp's sector  $H_\sigma$  after the change of variables  $\alpha_{\sigma(i)} = \prod_{j=i}^l \beta_j$  is performed,  $P$  factorizes as  $\prod_i \beta^{\omega(\sigma,i)} (a_\sigma + Q_\sigma(\beta))$  where  $a_\sigma$  is a non zero constant, and  $Q_\sigma$  is a polynomial in  $\beta$ 's. It is not too hard to check from (I.4.20) that  $U_G$  is fine; one identifies a particular spanning tree  $T_\sigma$  in  $G$  whose lines have maximal numbering with respect to the permutation  $\sigma$ ; the corresponding monomial gives the leading term of the fine property, hence in this case  $a_\sigma = 1$ ; then one checks that all other trees in  $G$  give non trivial monomials in  $\beta$ , hence their sum is the polynomial  $Q_\sigma$ .

With these indications the reader can easily recover Weinberg's theorem. But it is harder to give a bound on the value of the corresponding Feynman amplitudes which is of the type  $K^n$  where  $K$  is a constant and  $n$  the order of the graph (a good measure of its “size”). Such a bound, which we call the “uniform” Weinberg theorem, was in fact first proved in the  $\alpha$  representation [dCR] but this first proof is admittedly not very transparent. The goal of this section is to give in detail a more transparent proof, based on what we call a multiscale representation.

Multiscale decomposition inspired early work of K. Wilson on the renormalization group; introduced in constructive field theory by Glimm and Jaffe [GJ1], it was in particular developed and applied in [MS1-3][FO][Fed1]; the version we use now evolved from these works and follows [FMRS1] with some simplifications and improvements.

We consider a sequence of “momentum slices” which follow a geometric progression of ratio  $M$ , where  $M > 1$  is fixed (later it might be often convenient to choose  $M$  integer). If we use an ultraviolet cutoff in  $\alpha$  space of the type (I.3.7-8), it is natural to cut these slices in  $\alpha$  parametric space. So we start from (I.3.7), with  $\kappa = M^{-2\rho}$ . The propagator  $C_\kappa$ , now written  $C_\rho$  with a slight abuse of notation is cut into a discrete sum of propagators  $C^i$  by the rule:

$$C_\rho = \sum_{i=0}^{\rho} C^i \quad (\text{II.1.2})$$

$$C^i = \int_{M^{-2i}}^{M^{-2(i-1)}} e^{-m^2\alpha - \frac{|x-y|^2}{4\alpha}} \frac{d\alpha}{\alpha^{d/2}} \quad (\text{II.1.3})$$

$$C^0 = \int_1^\infty e^{-m^2\alpha - \frac{|x-y|^2}{4\alpha}} \frac{d\alpha}{\alpha^{d/2}}. \quad (\text{II.1.4})$$

$\alpha$  being dual to  $p^2$ , one should consider each propagator  $C^i$  as corresponding to a theory with both an ultraviolet and an infrared cutoff, which differ by the fixed multiplicative constant  $M$ , the momentum slice “thickness”.

The decomposition (II.1.2-4) is the fundamental technical tool used in this book. By the general structure of gaussian measures, there are associated decompositions of the gaussian measure  $d\mu_\rho$  corresponding to the propagator  $C_\rho$  into a product of independent gaussian measures  $d\mu^i$  and of the field  $\phi_\rho$ , as a random variable distributed according to  $d\mu_\rho$ , into a sum of independent random variables  $\phi^i$ , each  $\phi^i$  distributed according to  $d\mu^i$ :

$$\phi_\rho = \sum_{i=0}^{\rho} \phi^i; \quad d\mu_\rho(\phi_\rho) = \otimes_{i=0}^{\rho} d\mu^i(\phi^i) \quad (\text{II.1.5})$$

There are of course other technical ways to introduce momentum slices; in statistical mechanics or for lattice regularized models, the technique of block spin transformations [Ka] is very general and powerful; in constructive theory it has been used for instance in [BCGNOPS][GK1-4][Ba2-9][Fed2-7] etc...; it consists in writing, in a sequence of scaled lattices, each field variable as an averaged field in some cube of the next scale, plus a fluctuation field. An other elegant method close to (II.1.2-5) consists in defining the sliced fields as the coefficients of the full field on an orthonormal ondelettes basis [Bat2-3]. Each method has some advantages and drawbacks; the decomposition (II.1.2-5) is not the most general since it is only adapted to expansions around a gaussian measure, a frequent but not universal situation. Nevertheless it is simple, and has over the more general block spin methods the advantage of independence of the gaussian measures in various slices. We summarize this property by saying that the covariance is diagonal (in the space of values  $i$ ); this leads to a clear picture of the respective rôle of the gaussian measure and of the interaction, as will be stressed below.

In this part II devoted to perturbation theory, the language of fields and gaussian measures may be most of the time avoided; in Feynman diagrams the functional integral over the fields is already performed and only propagators do appear, so that we will use mostly (II.1.2-4). However it is a good thing to keep also in mind the related decomposition (II.1.5), since the field point of view becomes more and more important as harder problems are considered, and is absolutely necessary for constructive theory (part III).

It is an easy exercise to derive the following bound:

**Lemma II.1.1.**

There are positive constants\*  $K > 1$  and  $\delta < 1$  such that:

$$C^i(x, y) \leq KM^{(d-2)i}e^{-\delta M^i|x-y|} \quad (\text{II.1.6})$$

This bound captures the significant aspect of both cutoffs; an overall factor which shows that the singularity of  $C$  at coinciding points has been smoothed by the ultraviolet cutoff at a certain scale, and a scaled spatial decrease which comes from the infrared cutoff. (In the case  $i = 0$ ,  $\delta$  can be taken as any number less than  $m$ , the mass appearing in  $C$ ). This is optimal from the point of view of Fourier analysis; better “spatial resolution” costs a worse overall power counting factor.

Of course changing by a fixed factor  $M$  or  $M^{-1}$  the values of  $K$  and  $\delta$ , we could as well rewrite (II.1.6) as:

$$C^i(x, y) \leq KM^{(d-2)(i+1)}e^{-\delta M^{i+1}|x-y|} \quad (\text{II.1.7})$$

which is later convenient because  $i + 1$  is exactly the number of slices between 0 and  $i$ .

Using the slice decomposition we rewrite the bare amplitude for a Feynman graph as:

$$A_G = \sum_{\mu \in \mathbb{N}^{l(G)}} A_{G,\mu} \\ A_{G,\mu} = \int \prod_v dx_v \prod_{l \text{ internal line of } G} C^{i_l}(x_l, y_l) \prod_{l \text{ external line of } G} C_\rho(x_l, y_l) \quad (\text{II.1.8})$$

where  $\mu$  is called a “momentum assignment” (or “index assignment”) and is a list of integers, which gives for each internal line  $l$  of  $G$  the index  $i_l(\mu)$  (or simply  $i_l$ ) of the corresponding slice.  $A_{G,\mu}$  is the amplitude associated to the pair  $(G, \mu)$ , and (II.1.8) is called the multiscale representation of Feynman amplitudes. For the moment it is convenient to consider that external lines have an associated fictitious index -1, but that their propagator  $C_\rho$  is not decomposed into slices. Decomposition of external lines into slices may become useful for detailed results,

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\* **Important convention** We use  $\delta, \delta', \epsilon, \zeta, \dots$  as generic names for small constants and  $K, K_1, c, \dots$  for large ones. We may also write *const.* for a constant in some equations. We try to avoid any ambiguity, but these “constants” are often constant only in the context of the statement, and may depend in fact of the value of some parameters in a larger context.

e.g. on asymptotic behavior as some set of external momenta are scaled, but we do not use it in this book since we focus about basic existence problems.

The bare amplitude in (II.1.8) appears as a sum over the position of vertices and over momentum assignments. Together, these two sums make a kind of sum over “phase space”; but this traditional name is misleading, because *this* phase space is not the  $2d$ -dimensional cotangent bundle associated to  $d$ -dimensional space! It is rather a sort of  $d + 1$ -dimensional space with  $d$  continuous spatial dimensions and one discrete dimension, called the index space (or momentum space, with some abuse of language). This discrete space is moreover in most problems only a half-space, because we are usually interested either in an ultraviolet or in an infrared limit, but not in both together.

A key fact to realize is that in this phase space lines and vertices play a dual rôle. To support this intuition we draw a picture, Fig.II.1.2, of this phase space, which is our fundamental way of thinking to the decomposition (II.1.8); so we urge the reader to get some familiarity with it before going on, since we use the underlying intuition all the time in the rest of this book.

In the two dimensional plane of the figure, we use the horizontal direction to picture the  $d$  dimensions of space-time, and the vertical one to picture the discrete momentum slices, with the highest ones at the top. Then a propagator belongs to the slice of its index and appears as an horizontal line joining two vertices. Internal vertices sit at a particular point in space and join 4 half-lines which may be located in different slices. Hence it is convenient to picture them as vertical lines connecting the 4 horizontal half-lines hooked to them. These lines are dotted to distinguish them from the first ones. Finally the external lines are pictured in the fictitious “-1” slice, hence at the bottom of the picture.

Our goal is now to obtain the following theorem:

**Theorem II.1.1: Weinberg uniform theorem for  $\phi_4^4$**

There exists a constant  $K$  such that for any connected completely convergent  $\phi_4^4$  graph (i.e.  $\omega(S) > 0$ ,  $\forall S \subset G$ ) the Feynman amplitude (I.4.8) of  $G$  is absolutely convergent and the following bound holds:

$$|A_G| \leq K^{n(G)} \cdot \text{Ext} \quad (\text{II.1.9})$$

where the function  $\text{Ext}$  depends on the way external arguments are treated.

This theorem extends to a very general class of completely convergent graphs which have power counting properties of the just renormalizable kind, but are not necessarily  $\phi^4$  graphs. In this more general version, in which an unbounded number of lines may be hooked to a single vertex, it is  $l(G)$ , not  $n(G)$  which is the natural measure of the size of the graph, so that the corresponding generalized bound is similar to (II.1.9) but with a factor  $K^{l(G)}$  instead of  $K^{n(G)}$ ; by formula (I.4.3), both bounds are equivalent for  $\phi^4$ . For this more general result we refer to [FMRS1].

By factorization, bounds for non-connected graphs can also be derived easily from (II.1.9). The point to emphasize in (II.1.9) is the uniform exponential character of the bound at large order expressed by the factor  $K^{n(G)}$ ; the particular form of  $\text{Ext}$  is secondary. In [FMRS1] the set of external vertices is called  $V_E$ , and

three possible treatments for external cases are considered (the reader is invited to try its own favorite):

- H1 Each external vertex is integrated over a standard unit cube of  $\mathbb{R}^4$ ,
- H2 Each external vertex  $v$  is integrated against a test function  $f_v$  on  $\mathbb{R}^4$ ,
- H3 Each external vertex has a fixed external momentum  $k_v$  entering it.

In these three cases a possible choice for  $Ext$  would be:

- H1  $Ext = \sup_{x_v, v \in V_E} \prod_{l \in G} e^{-m(1-\zeta)|x_l - y_l|}$
- H2  $Ext = \inf_{v \in V_E} \|f_v\|_{L^1} \prod_{w \neq v \in V_E} \|f_w\|_{L^\infty}$
- H3  $Ext = \delta(\sum_{v \in V_E} k_v)$

To underline the essential part of the argument, let us prove the theorem in the slightly simpler case of an amplitude for which external propagators have been amputated, and exactly one internal vertex  $v_0$  is fixed to the origin (not integrated over  $\mathbb{R}^4$ ). This is in essence the case H3, because each external propagator  $(p^2 + m^2)^{-1}$  is bounded by  $m^{-2}$  and  $N(G) \leq 4n(G)$ , so that amputation does not change the nature of the bound at large order, and fixing  $v_0$  at the origin is the same as taking into account the overall  $\delta$  function of  $Ext$ , which simply reflects translation invariance.

We consider (II.1.8) and perform first the integration over the positions of the vertices in  $\mathbb{R}^4$ . To do this in the best possible way according to the momentum assignment, one should use as much as possible the decay of the lines with highest possible index (see (II.1.6-7)). This leads naturally to consider, for a given  $\mu$  and  $i \in \mathbb{N}$ , the connected components of  $G^i$ , the subgraph of  $G$  made of all lines with index  $j \geq i$ . Let us call these connected components  $G_k^i$ ,  $k = 1, \dots, l(i)$ . There is a systematic way to know whether a given connected subgraph  $g \subseteq G$  is a  $G_k^i$  for some  $i$  and  $k$ . We define the internal and external index for  $g$  in the assignment  $\mu$  as:

$$i_g(\mu) = \inf_{l \in g} i_l(\mu) \quad (\text{II.1.10})$$

$$e_g(\mu) = \sup_{l \text{ external line of } g} i_l(\mu) \quad (\text{II.1.11})$$

(with the  $\mu$  dependence sometimes omitted for shortness). The condition that in the assignment  $\mu$  the subgraph  $g$  is a  $G_k^i$  for some values of  $i$  and  $k$  is simply:

$$i_g(\mu) > e_g(\mu) \quad (\text{almost local condition}) \quad (\text{II.1.12})$$

Subgraphs verifying condition (II.1.12) are called *almost local* (with respect to the assignment  $\mu$ ). For such a  $g$  and each value of  $i$  such that  $e_g(\mu) < i \leq i_g(\mu)$  there exists a value of  $k$  such that  $g$  is equal to  $G_k^i$ .

The  $G_k^i$ 's are partially ordered by inclusion and form in fact a forest, i.e. a set of connected graphs such that any two of them are either disjoint or included one into the other (disjoint here means not only no line but also no vertex in common). This remark is essential both here and for the next section, which deals with renormalization. Since  $G = G^0$  is itself connected, the forest is in fact a tree, whose root is the full graph  $G$ . This tree as been pictured in Fig.II.1.3 for the case corresponding to  $G$  and  $\mu$  as in Fig. II.1.2.

Each  $G_k^i$  is a node on this picture and a line between  $G_k^i$  and  $G_{k'}^{i-1}$  simply means that  $G_k^i \subseteq G_{k'}^{i-1}$ . For  $g$  almost local there is a sequence of vertical lines in this tree with  $i_g - e_g$  nodes corresponding to  $g$ . This sequence may be collapsed to a single line for simplicity and we obtain the tree of almost local subgraphs pictured in Fig.II.1.3.

The tree of Fig.II.1.3 is not exactly the same but closely related to the Gallavotti-Nicoló trees which are the key tools in their analysis of renormalization theory [GaNi][Gal].

To integrate over the positions of internal vertices (save one,  $v_0$ ) requires at least the decay of a spanning tree, which is a minimal set of lines connecting together all the vertices. For such a spanning tree  $T$ , defining the  $x_0$  vertex to be the root induces a particular partial ordering on the vertices of  $G$ , as explained in Sect.I.4. Such a particular tree for the situation of Fig.II.1.2-3 is pictured in Fig.II.1.4.

Remark that the two trees of Fig.II.1.3 and Fig.II.1.4. have very different meaning; the one in Fig.II.1.3 is an abstract picture of the inclusion relations derived from assignment  $\mu$ , but the tree of Fig.II.1.4 is a concrete set of lines of  $G$ . Let us call  $i_\sigma$  the index in  $\mu$  of the line directly under vertex  $v$  in the partial ordering shown in Fig.II.1.4. If we use the decay of the lines of  $T$  to integrate over the positions of the vertices, starting from the “leaves” at the top of the tree towards the root which is fixed, we have to forget the decay of the other lines of  $G$  (the ones not in  $T$ ) in order to have a simple expression. This is possible for an upper bound (but would be difficult for a lower bound). This means that for any  $T$  we can write:

$$\begin{aligned} \int \prod_{v \neq v_0} d^4 x_v \prod_{l \in G} e^{-\delta \sum M^{i_l+1} |x_l - y_l|} &\leq \int \prod_{v \neq v_0} d^4 x_v \prod_{l \in T} e^{-\delta \sum M^{i_l+1} |x_l - y_l|} \\ &= \prod_{v \neq v_0} K M^{-4(i_\sigma+1)} \end{aligned} \quad (\text{II.1.13})$$

where  $K$  is a constant (here  $\delta^{-4} \cdot \Omega_3$  where  $\Omega_3$  is the volume of the unit three dimensional sphere).

Now we should choose  $T$  so as to optimize (II.1.13), i.e. to make the  $i_\sigma$ 's as large as possible. The best possible choice would be to require that  $T$  restricted to every almost local subgraph is still a spanning tree of this subgraph. In this case there is a kind of maximal compatibility between the trees of Fig.II.1.3 and II.1.4. Such an optimal choice is always possible (usually in a non unique way). Indeed we can first choose spanning trees  $T_k^\rho$  in each “leaf”  $G_k^\rho$  where  $\rho$  is the maximal index appearing in  $\mu$ ; then by induction, when several  $G_{k'}^i$  are included in a single  $G_k^{i-1}$ , we can always add one line in  $G_k^{i-1}$  to the union of the corresponding  $T_{k'}^i$ 's without creating loops in  $G_k^{i-1}$ ; we can repeat this until it is impossible to add lines without creating a closed loop; but then we must have built a spanning tree  $T_k^{i-1}$  of  $G_k^{i-1}$ . In this way a global  $T$  is built inductively, which has the required property of spanning each  $G_k^i$ . Remark that the tree structure of the  $G_k^i$  themselves is crucial for this optimization to be possible; obviously for the bubble graph of Fig.I.1.1 there is no way to choose a tree whose restriction to the two “overlapping” subsets made of one line of  $G$  is a spanning tree of both.

Once this (non unique) choice of  $T$  is completed, for any  $G_k^i$  every vertex  $v$  save one has  $i_v \geq i$ . (This is an easy consequence of the connectedness of  $T \cap G_k^i$ ). We use (II.1.8) and (II.1.13) with this choice of  $T$ . Apart from overall constants of type  $K^{n(G)}$  or  $K^{l(G)}$  (which are allowed in the bound since  $l(G) \leq 2n(G)$ ), the contribution  $A_{G,\mu}$  is bounded by:

$$\prod_{i \geq 0} \prod_{k=1}^{l(i)} \left\{ \left[ \prod_{l \in G_k^i} M^2 \right] M^{-4[n(G_k^i)-1]} \right\} = \prod_{i \geq 0} \prod_{k=1}^{l(i)} M^{-\omega(G_k^i)} \quad (\text{II.1.14})$$

To get this bound one should first rewrite for each line  $l$ :

$$M^{2(i_l+1)} = \prod_{(i,k) \text{ such that } l \in G_k^i} M^2$$

because for every  $i \leq i_l$ , the line  $l$  belongs to exactly one  $G_k^i$ , and to none for  $i > i_l$ .

Now recall that since  $G$  is completely convergent we have:

$$\omega(G_k^i) = N(G_k^i) - 4 \geq \frac{N(G_k^i)}{3} \quad (\text{II.1.15})$$

Let us define

$$e_v(\mu) = \sup_{l \text{ hooked to } v} i_l(\mu) \quad (\text{II.1.16})$$

$$i_v(\mu) = \inf_{l \text{ hooked to } v} i_l(\mu) \quad (\text{II.1.17})$$

where the inf in (II.1.17) is over *every* line hooked to  $v$ , including external lines, which by convention have index  $-1$ .

We remark that for any  $i$ , a given vertex  $v$  belongs to exactly one  $G_k^i$  for  $i \leq e_v(\mu)$  and to none otherwise. Furthermore some external lines of this  $G_k^i$  are hooked precisely at  $v$  if and only if  $i_v(\mu) < i \leq e_v(\mu)$ . Hence, using (II.1.15):

$$\prod_{i \geq 0} \prod_{k=1}^{l(i)} \{M^{-\omega(G_k^i)}\} \leq \prod_{i \geq 0} \prod_{k=1}^{l(i)} \{M^{-N(G_k^i)/3}\} \leq \prod_v M^{-\frac{1}{3}|e_v(\mu) - i_v(\mu)|} \quad (\text{II.1.18})$$

The meaning of this bound is that for a completely convergent graph, i.e. one which has favorable power counting, after spatial integration the vertices pictured as dotted lines in Fig.II.1.2 acquire an exponential decay in their length  $|e_v(\mu) - i_v(\mu)|$  in the vertical direction. This is in a sense simply power counting viewed on a logarithmic momentum scale (the space of discrete indices of the slices), but a power counting ready to be exploited. Therefore the main advantage of the multiscale decomposition emerges: it may be viewed as a machine which at each scale effectuates some spatial integration affected to this scale and, if power counting is favorable, gives as a reward a small factor for each dotted line (vertex) which crosses this scale. These small factors build up a vertical exponential decay which is very much the dual of the horizontal exponential decay of the ordinary lines in each slice. It is intuitively obvious that this decay should make the sum over momentum assignments easy, the external lines with their index  $-1$  breaking

vertical translation invariance, and playing therefore a rôle dual to the one of the fixed vertex  $x_0$ .

There are several ways to make this intuition more precise; let us describe one of them, with no efforts to find optimal constants. Using the fact that there are at most 4 half-lines, hence at most 6 pairs of half-lines hooked to a given vertex, and that for such a pair obviously  $|e_v(\mu) - i_v(\mu)| \geq |i_l - i_{l'}|$ , we can convert the decay in vertical length of (II.1.18) into a decay associated to each pair of half-lines hooked to the same vertex:

$$\prod_v M^{-\frac{1}{3}|e_v(\mu) - i_v(\mu)|} \leq \prod_v \left\{ \prod_{(l, l') \text{ hooked to } v} M^{-\frac{1}{18}|i_l - i_{l'}|} \right\} \quad (\text{II.1.19})$$

where again lines  $l, l'$  in (II.1.19) are either internal or external, and the factor 18 is not optimal (it can be improved to 12 with negligible effort...). The analog of picking a tree to perform the spatial integration is to pick a total ordering of the internal lines of  $G$  as  $l_1, \dots, l_{l(G)}$  such that  $l_1$  is hooked to an external line of  $G$  and such that each subset  $\{l_1, \dots, l_m\}$ ,  $m \leq l(G)$  is connected, which is clearly possible. Using only a fraction of the decay in (II.1.19) we have:

$$\prod_v \left\{ \prod_{(l, l') \text{ hooked to } v} M^{-\frac{1}{18}|i_l - i_{l'}|} \right\} \leq \prod_{j=1}^{l(G)} M^{-\frac{1}{18}|i_{l_j} - i_{l'_{j'(j)}}|} \quad (\text{II.1.20})$$

where  $j'(j)$  is the index of a line hooked at  $j$  but of lower index:  $1 \leq j'(j) < j$  if  $2 \leq j \leq l(G)$  and by convention  $i_{l'_{j'(1)}} = -1$ .

Finally in (II.1.20) performing the sum over  $\mu$  we obtain:

$$A_G \leq K_1^{n(G)} \sum_{\mu=\{i_1, \dots, i_{l(G)}\}} M^{-\frac{1}{18}|i_{l_j} - i_{l'_{j'(j)}}|} \leq K_1^{n(G)} K_2^{l(G)} \leq K^{n(G)} \quad (\text{II.1.21})$$

which achieves the proof of the theorem.

This proof is elementary and quite compact, but a subtle point remains, which is the optimization over the trees  $T$ . It introduces some sort of global choice which is not very transparent. Therefore without giving all details we describe now another slightly different way of getting the same result, which although in a sense more complicated, has the advantage of being a single induction which works from higher frequencies towards the lower ones. One can at each scale  $i$  establish an inductive bound in which spatial integration is completed for positions of all vertices save one in each  $G_k^i$ , *without knowing the structure at lower scales*. The general outline is as above: In each  $G_k^i$  all vertices save one should be integrated with lines of scales higher or equal to  $i$ , and this gives a volume effect naively in  $M^{-4i \sum_k (n(G_k^i) - 1)}$ , but when powers of  $M$  are distributed over all scales as above, it is rather  $M^{-4 \sum_k (n(G_k^i) - 1)}$  which should be attributed to scale  $i$ . Similarly out of the  $M^{2i_l}$  associated to a line  $l$  should come for scale  $i$  a factor  $M^{+2 \sum_k l(G_k^i)}$ ; together they reconstruct the desired factor  $M^{-\sum_k \omega(G_k^i)}$ , from which vertical exponential decrease for dotted vertical lines and a uniform bound on the sum over momentum assignments follows again. However the difficulty in this point

of view in which no global choice of  $T$  is made first, is that one should preserve a piece of the internal decay between the vertices of  $G_k^i$  for later use, otherwise spatial integrations over vertices in some  $G_{k'}^{i'}$ , with  $i' < i$  may be blocked later (if these vertices were joined in  $G_{k'}^{i'}$  only by paths passing through lines of  $G_k^i$ , which is perfectly possible). But how to remember a decay when points are supposed to have been already integrated out? It is of course possible but there is no canonical way. In the case of the exponential decay (II.1.6) or (II.1.7) a possible technical solution is to duplicate the decrease of each line  $l$  into a product over all scales smaller or equal to  $i_l$ , since a geometric sum is of same order as its leading term. Hence changing  $\delta$  one writes:

$$C^i(x, y) \leq KM^{(d-2)(i+1)} e^{-\delta \sum_{0 \leq j \leq i} M^{j+1} |x-y|} \quad (\text{II.1.22})$$

Then as integration is made with the help of the higher scale decay, a simple supremum may be taken over the lower scales decay. This supremum becomes available later to reconstruct tree decay of the right scale between the remaining vertices of the  $G_{k'}^{i'}$  even after the ones of  $G_k^i$  (save one) have been integrated [FMRS1].

In spite of this technical complication, the inductive point of view is important because it adapts better to constructive field theory and to the general philosophy of the renormalization group, in which high frequencies should be integrated out without information about lower ones.

One should emphasize that the exponential spatial decay (II.1.6) is not crucial to the phase space machinery. In particular a scaled power decay of the type:

$$C^i(x - y) \leq KM^{2i} [1 + M^i |x - y|]^{-p} \quad (\text{II.1.23})$$

is sufficient provided  $p > 4$  (summability of the propagator is what really matters), and it becomes almost as convenient as exponential decay if  $p$  is very large. Different type of decay may occur if different rules for the momentum slicing are used; in particular one may prefer to slice with a  $C_0^\infty$  cutoff in momentum space, because one can then bound with probability one norms on the field which contain derivatives, like Sobolev norms, by ordinary norms like  $L^2$  norms. For such cutoffs exponential decay in  $x$ -space is no longer possible, but power decay like (II.1.23) can still hold. An interesting exercise suggested to the reader is to rewrite in this case the inductive proof sketched above; the replica trick in the simple form of (II.1.22) is no longer allowed!

Finally it is interesting to compare this multiscale proof with earlier proofs of similar results by constructivists for  $\phi_3^4$  [Gl]. Here like in any superrenormalizable theory one can exploit the fact that after some minimum size is attained a graph has to be superficially convergent. But amplitudes for subgraphs can be viewed as functions of the positions of their “border vertices” or as kernels for integral operators. In superrenormalizable theories one can devise a way to break a big graph into a lot of pieces of finite size, large enough so that the kernel for each piece has a finite Hilbert-Schmidt norm. Using convex analysis for instance in the form of some Schwarz inequalities a uniform bound similar to (II.1.9) can be reached, the constant  $K$  being related to a supremum over the *finite* set of these

typical pieces. This method and its extension to functional integrals was widely used in the first period of constructive theory.

But no refinement of this method is going to work for a theory like  $\phi_4^4$ , because Schwarz or other convex inequalities cannot be applied without losing a fraction of the power counting, and this loss is fatal to a just renormalizable theory. Multiscale chopping becomes the right tool, and once adopted it seems simpler to apply it to the superrenormalizable case as well. We leave as an exercise to the reader to check that in this superrenormalizable case the vertical decrease occurs for every vertex from  $e_v(\mu)$  all the way down to  $i = 0$ ; hence, as long as their decay is concerned, dotted lines of decay may be drawn in this case from the highest line at a vertex down to the bottom (the slice with  $i = 0$ ). This result will be used in section III.1D to construct the ultraviolet limit of the  $\phi_2^4$  theory in a unit square.

Comparing the dotted lines to some sort of springs, the kind of intuition to be gained from this exercise and the rest of this section is that for completely convergent graphs, in the superrenormalizable case any subgraph is tied to the “ground” of low energies (the  $i = 0$  slice) by springs attached to all its vertices, whereas in the just renormalizable case it is tied only by springs at its “border vertices”, and not directly to the ground; chains of such springs will ultimately reach the ground at some true external lines of the full graph but the itinerary for that may be arbitrarily long and complex, involving many intermediary other subgraphs. Nevertheless in both cases the same uniform bound ultimately holds. This intuitive picture is also very useful for the next section in which multiscale slicing will lead us to an effective analysis of renormalization.

## II.2 Renormalization theory for $\phi_4^4$

We start with an overview of the situation and some examples before to introduce the full formalism of renormalization, which unfortunately even in its most recent and transparent versions still involves some heavy notations.

Our first remark is that the proof we gave of the “uniform Weinberg” theorem also gives uniform bounds for many momentum assignments corresponding to graphs which are no longer completely convergent. Indeed the only convergence degrees which appear in a given  $\mu$  are the  $\omega(G_k^i)$ ’s. This leads us to call a momentum assignment  $\mu$  for a graph  $G$  *convergent* if  $\omega(g) > 0$  for any  $g$  almost local with respect to  $\mu$ . Then in the preceding section we have proved in fact:

$$|A_G^c| = \left| \sum_{\mu \text{ convergent}} A_{G,\mu} \right| \leq K^{n(G)} \quad (\text{II.2.1})$$

Remark that this result, strictly speaking, is non trivial only for superficially convergent  $G$ ’s, otherwise there are no convergent assignments (since  $G$  itself is always almost local, by our convention that the external lines of  $G$  are considered to have index  $-1$ ).

When some  $G_k^i$  have  $N(G_k^i) = 4$  or  $2$  ( $0$  is excluded if  $G$  is connected and  $N(G) > 0$ ) the proof of the last section breaks down. In the intuitive language we used, the phase space machinery, instead of rewarding by a small factor for each corresponding pair  $i$  and  $k$  becomes neutral (if  $N(G_k^i) = 4$ ) or even costs a large factor (if  $N(G_k^i) = 2$ ). In the language of the renormalization group these behaviors are called respectively marginal and relevant; they are the source of the famous ultraviolet divergences.

What does the phase space point of view teach us about these divergences? There are two important facts to remark, which pave the way for the full solution of the problem of perturbative renormalization, and both facts appear particularly clearly in the multiscale representation.

First the divergences arise from the appearance of superficially divergent (or “divergent”, for short) almost local subgraphs  $G_k^i$ ’s. Their name has been chosen to suggest that these objects are almost point-like from the “external world” point of view. By this we mean that their internal lines, because of the decay (II.1.6) may be thought as horizontal springs which constrain them to extend only over a region in space typically of diameter of the order of  $M^{-i}$ . In contrast the external lines which in an intuitive sense carry out information about these subgraphs to the rest of the graph, are in a lower horizontal slice, hence they distinguish only larger scales. So for them the  $G_k^i$ ’s appear *almost* point-like, and this effect becomes stronger as the gap grows between the internal and external scales which is precisely the source of ultraviolet divergences. It is therefore not surprising that these divergences can be cancelled by comparing these contributions to the ones of a purely local (*exactly* point-like) counterterm, which is precisely what renormalization does.

The second and more subtle remark is that the divergent almost local subgraphs have obviously a forest structure inherited from the tree structure of the inclusion relations of all  $G_k^i$ ’s. This simple remark is the key to the solution of one of the historically most subtle and debated point in the theory of perturbative renormalization, the one of “overlapping divergences”. It was early understood

and later formalized by Bogoliubov and others [BP][BS][He1][He2][Zim] that every renormalization implemented by counterterms in the Lagrangian gives rise to contributions in the amplitudes indexed by forests of connected divergent subgraphs. This is not really difficult to grasp. Indeed local counterterms in the lagrangian are like vertices; branching propagators on them will create contributions which can match any structure of divergent disjoint subgraphs (see Fig.II.2.1). Since the definition of counterterms may be inductive one can imagine that their definition includes earlier subtractions inside them, and in this way contributions may be associated with stacks of subgraphs included one into the other as in Fig.II.2.2. Combining both ideas one gets contributions for each forest. But there is certainly no way to associate a specific contribution which in the perturbative expansion matches with two divergent overlapping subgraphs (i.e. with nontrivial intersection, see Fig.II.2.3.). For many years experts worried about these overlapping divergences, since there was no way to built a specific contribution for them from local counterterms. The fact that the almost local subgraphs have forest structure is precisely the solution to this point; “overlapping divergences” simply never occur simultaneously in phase space, hence there is no need to associate any specific counterterm to them. This statement will be fully substantiated below.

After this general overview, we start by showing how the divergent “bubble” graph of Fig.I.1.1 is renormalized by a local counterterm in the multiscale representation. A subgraph  $g \subset G$ , after integration of its inside vertices (it may have none) becomes a function  $g(x_1, \dots, x_{v_e})$  of its  $v_e$  border vertices (see Sect.I.4 for these notions). By translation invariance its Fourier transform is of the form  $\delta(k_1 + \dots + k_{v_e})\hat{g}(k_1 + \dots + k_{v_e-1})$ . The Zimmermann prescription for renormalizing it at 0 external momenta is to subtract from it a counterterm  $\tau_g g$  which is  $\delta(k_1 + \dots + k_{v_e})$  times a Taylor expansion of  $\hat{g}$  around 0 momenta:

$$\sum_{j=0}^{D(g)} \frac{1}{j!} \frac{d^j}{dt^j} \hat{g}(tk_1, \dots, tk_{v_e-1})|_{t=0} \quad (\text{II.2.2})$$

where  $D(g)$  is the largest integer less than or equal to  $-\omega(g)$  the superficial degree of *divergence* of  $g$ . Thanks to translation invariance, the formula does not really depend on  $v_e$ . Remark that the 0-momentum subtraction scheme is valid only for a massive theory; otherwise it would be ill defined, due to infrared singularities, and some subtraction scale should be introduced by hand to represent the scale of physical phenomena one is interested in. This is for instance the case for non-abelian gauge theories, in which the gauge boson is massless. (A subtraction scheme for massless  $\phi_4^4$  and corresponding bounds may be found in [dCPR]).

In the multiscale representation we want to compute in  $x$  space; therefore we should apply a subtraction operator  $\tau_g^*$  to the test function smearing  $g$  rather than to  $g$  itself. This is possible through the formula:

$$\begin{aligned} & \int \tau_g g(x_1, \dots, x_{v_e}) a(x_1, \dots, x_{v_e}) dx_1 \dots dx_{v_e} \\ &= \int g(x_1, \dots, x_{v_e}) \tau_g^*(v_e) a(x_1, \dots, x_{v_e}) dx_1 \dots dx_{v_e} \end{aligned} \quad (\text{II.2.3})$$

where

$$\tau_g^*(v_e)a(x_1, \dots, x_{v_e}) = \sum_{j=0}^{D(g)} \frac{1}{j!} \frac{d^j}{dt^j} a(x_1(t), \dots, x_{v_e}(t))|_{t=0} \quad (\text{II.2.4})$$

with  $x_i(t) = x_{v_e} + t(x_i - x_{v_e})$ . Here  $v_e$  really plays a distinguished rôle, so that there are several noncanonical ways of defining the “adjoint”  $\tau_g^*(v_e)$ , one for each border vertex  $v_e$ ; they are equally good.

Let us consider a bubble subgraph  $g$  in a bigger graph  $G$ . In this case  $\omega(g) = 0$  and there are two border vertices at  $x_1$  and  $x_2$ . Let us suppose to simplify that the two internal lines are propagators of the same scale  $i$  and that the four external half-legs are four different propagators of the same scale  $j$ . The test function  $a$  is therefore the product  $C^j(x_1, y_1)C^j(x_1, y_2)C^j(x_2, y_3)C^j(x_2, y_4)$  and the subgraph is almost local (divergent case) iff  $j < i$  (see Fig.II.2.4).

As we saw in the last section there is no decay in  $i-j$  in this case, which would allow to sum over the  $i$  index, holding  $j$  fixed, hence there is a “logarithmic” divergence in energy. To renormalize, we apply  $(1 - \tau_g^*)$  to  $a$ , which gives the bare amplitude minus its counterterm. The result may be written as the following sum:

$$\begin{aligned} & [C^j(x_1, y_1) - C^j(x_2, y_1)]C^j(x_1, y_2)C^j(x_2, y_3)C^j(x_2, y_4) \\ & + C^j(x_2, y_1)[C^j(x_1, y_2) - C^j(x_2, y_2)]C^j(x_2, y_3)C^j(x_2, y_4) \end{aligned} \quad (\text{II.2.5})$$

Each term contains exactly one difference which may be written as a Taylor remainder. For instance:

$$C^j(x_2, y_1) - C^j(x_1, y_1) = \int_0^1 dt \frac{d}{dt} C^j(x_1 + t(x_2 - x_1), y_1) \quad (\text{II.2.6})$$

Returning to the definition (II.1.3-4) of  $C^j$  it is easy to derive the bound:

$$\frac{d}{dt} C^j(x_1 + t(x_2 - x_1), y_1) \leq K|x_2 - x_1|M^{3j}e^{-\delta M^j|x_1+t(x_2-x_1)-y_1|} \quad (\text{II.2.7})$$

where  $K$  and  $\delta$  are new constants, which may be taken slightly worse than those of (II.1.6). Hence the net effect of the derivation is to add to the usual bound a factor  $M^j$  times the distance  $|x_2 - x_1|$ . This distance, between the two border vertices of the bubble may be then estimated by using a piece of the decay of the two internal lines of the bubble. Remember that these lines have decay of scale  $i > j$ . Hence this gives a bound:

$$|x_2 - x_1|e^{-2\delta M^i|x_2 - x_1|} \leq \frac{2}{\delta} M^{-i} e^{-(3/2)\delta M^i|x_2 - x_1|} \quad (\text{II.2.8})$$

Similarly this internal decay can be used to replace  $|x_1 + t(x_2 - x_1) - y_1|$  in (II.2.7) by  $|x_1 - y_1|$ ; then the integral over  $t$  in (II.2.6) is just bounded by 1. Altogether renormalization has delivered an extra factor  $M^{-(i-j)}$  at the cost of making constants like  $K$  or  $\delta$  slightly worse. This does not change the nature of the estimates, but the extra factor  $M^{-(i-j)}$  has restored a vertical exponential decrease previously missing; more concretely it allows now to sum over  $i$  with respect to  $j$ .

In the generic case recall that the condition for a divergent subgraph  $g$  to be almost local with respect to  $\mu$  hence to be a source of ultraviolet divergence, is  $i_g(\mu) > e_g(\mu)$ . For instance in the case of a bubble subgraph, the condition of almost locality is simply that  $i_g$ , the smallest of the two indices  $i_1$  and  $i_2$  in  $\mu$  of its internal lines, has to be larger than  $e_g$ , the largest of the 4 indices  $e_1, e_2, e_3$  and  $e_4$  in  $\mu$  of its external lines. For  $g$  a general divergent subgraph, assignments  $\mu$  satisfying the condition  $i_g(\mu) > e_g(\mu)$  are called *dangerous* (with respect to  $g$ ) and the ultraviolet divergence for  $g$  solely comes from such assignments. In the case of a single bubble subgraph, the reader can easily extend the previous argument to the sum over such dangerous assignments and restore the missing decay in  $i_g - e_g$  by subtracting again the counterterm  $\tau_g$ . In conclusion and anticipating a bit, the divergence of the bubble is cancelled successfully by a counterterm in the Lagrangian of the form:

$$\int dx_1 \left[ \int dx_2 \sum_{i_g > e_g} C^{i_1}(x_1 - x_2) C^{i_2}(x_1 - x_2) \phi_{e_1}(x_1) \phi_{e_2}(x_1) \phi_{e_3}(x_1) \phi_{e_4}(x_1) \right] \quad (\text{II.2.9})$$

where we recall that the random field  $\phi^i$  of frequency  $i$  (respectively  $\phi_i$ ) is distributed according to the gaussian measure of covariance  $C^i$  (respectively  $C_i = \sum_{j=0}^i C^j$ ) (see (II.1.5)). We refer to pieces of counterterms like the one of (II.2.9) loosely as to “useful” counterterms, because they really kill a divergence. But objects like (II.2.9) cannot be considered local because of the restriction in the sum, which makes the value of the counterterm dependent of the energy scale of its external lines. If we want to preserve the formal locality of the Lagrangian, we have to introduce the counterterm for  $g$  also in the case of assignments which are not dangerous for  $g$ , hence satisfy the opposite condition in the sum (II.2.9); we call this piece of the counterterm the “useless” counterterm (it is “useless” from the point of view of cancellation of divergences, not locality).

To combine useless counterterms with the bare amplitude in assignments which are not dangerous does not make sense from the point of view of estimates, essentially because subgraphs which are not  $G_k^i$ 's are not almost local from the external point of view, and there is no interesting gain to compare them with a completely local counterterm.

An interesting “naive” question is: for these assignments, where they do not match, which one is biggest, the bare amplitude or the counterterm? A good short answer is: “the counterterm”. To motivate this answer we remark that the condition to be “useless”,  $i_g \leq e_g$ , puts an ultraviolet cutoff at scale  $e_g$  on the divergent sum over internal scales in the counterterm. This cutoff makes this sum finite, but there is a subtle point: apart from the small breaking of scale invariance due to the mass  $m$  in the propagator  $C$ , any counterterm for a marginal operator (like  $\phi^4$  here) is invariant by vertical translation. Hence when  $e_g$  is large, the useless counterterm has to behave asymptotically as  $e_g$ . But there is no such effect for the corresponding piece of the bare amplitude! This may seem a bit mysterious, but the reader may convince himself of this important fact by checking again the proof of the last section and its extension to (II.2.1) above: no linear divergence in the summation over scales can occur for subgraphs which are not  $G_k^i$ 's, because such subgraphs simply never appear in the argument. We conclude that at least for  $e_g$  large the useless counterterm must be much larger than the corresponding

bare contribution.

It is a bit uncomfortable that something which is useless from the point of view of convergence dominates. It is not hard to guess more precisely to which problem this phenomenon will lead us. Useless counterterms will not make amplitudes divergent; they are big, but finite. But what will happen if a large number of such big objects are inserted into a convergent bare amplitude or an amplitude with a useful renormalization? The convergence we found for such amplitudes is exponential in the space of momentum indices, hence it cannot be destroyed by multiplying by any number of linear factors like those associated to the useless counterterms. This is in essence why renormalization works and renormalized amplitudes are finite. However since:

$$\sum_{e=0}^{\infty} M^{-\delta e} e^n \simeq K^n n! \quad (\text{II.2.10})$$

we should expect the production of factorials in the number of useless counterterms, and this leads to a violation of the uniform bound (II.2.1) for renormalized amplitudes. This violation, called the “renormalon” phenomenon, leads to new difficulties if one tries to sum up renormalized perturbation theory. We have identified the useless counterterms as the source of this trouble, and this diagnosis already suggests that the renormalized series are not the ones that one should try to sum up, a point of view which will be developed in Sect.II.4-5.

Having completed this brief sketch of the situation, based on the sole example of the bubble, we are ready to enter into the heart of renormalization to substantiate the corresponding ideas with some proofs. We will focus on a proof of finiteness of renormalized amplitudes with reasonable estimates which imply the existence of a finite disk of analyticity for the Borel transform of the renormalized series. This result, the “uniform BPH theorem” below, was first obtained in [dCR1] using the  $\alpha$  space representation. Later it was rederived in several different ways [GaNi][FMRS2][FHRW][Hu]. As remarked already, a version like [Hu] may be the most compact. But the one given here, in the style of [FMRS2] is well suited to the extension to constructive theory. Small details have been improved so that the final estimates are slightly better than the ones of [dCR1] or [FMRS2]

The general statement that  $\phi^4$  is perturbatively renormalizable of course does not mean simply that it is possible to make amplitudes finite by some subtraction process, since this would be obvious for any theory in any dimension. It means that the subtraction process must correspond solely to counterterms in the action of the form  $\phi^4$ ,  $\phi^2$  and  $\partial_\mu \phi \partial^\mu \phi$ , as the original ones in (I.3.1).

More precisely it means that in (I.3.1) we can replace the constants  $g$ ,  $m^2$  and  $a$  of the theory by three formal power series in a renormalized coupling constant  $g_r$ , respectively  $g_r + \sum_{n=2}^{\infty} c_n g_r^n$ ,  $m^2 + \sum_{n=1}^{\infty} d_n g_r^n$  and  $a + \sum_{n=2}^{\infty} e_n g_r^n$  such that the perturbative expansion in  $g_r$  of any Schwinger function is finite order by order. The three corresponding renormalizations are called respectively coupling constant, mass and wave function renormalization. (For connected amplitudes there is no need to discuss the vacuum energy renormalization.) But counterterms  $c_n$ ,  $d_n$  or  $e_n$  should diverge in the continuum limit, so that this is not yet a rigorous definition. What is finally required is that for the well defined theory with cutoffs

like  $(\Lambda, \kappa)$  (I.3.9) (or  $(\Lambda, \delta)$  (I.3.12)) there exist three such formal power series with coefficients  $c_n(\Lambda, \kappa)$ ,  $d_n(\Lambda, \kappa)$  and  $e_n(\Lambda, \kappa)$  such that the perturbative expansion for this theory with cutoffs has a finite limit, order by order in  $g_r$ , when both cutoffs are removed.

In the standard litterature, renormalizability consists therefore of two steps: the first one is the definition of the coefficients  $c_n$ ,  $d_n$  and  $e_n$ , at the heart of which lies the Bogoliubov recursion; the second step is to expand the full theory in powers of  $g_r$  and to match the contributions obtained from the counterterms in the action coming from this recursion with the bare contributions so that every renormalized amplitude remains finite when cutoffs are removed. At the heart of this second step is a technical tool, Zimmermann's forest formula, which gives for a given graph the complete list of counterterms which should be combined with it; the forest formula appears therefore as a kind of inverse solution to the Bogoliubov recursion, and it has the advantage of being global, not inductive. Both tools unfortunately were invented before the multiscale representation and do not take advantage of it in their design.

This defect is alleviated in the recent versions of the renormalizability theorem [GaNi] [Gal] [FHRW] [Hu]; they neither focus on the Bogoliubov formula and Zimmermann's forest, nor on the individual renormalized graphs. To the decomposition of the propagator into slices corresponds the decomposition of the field as a sum of random variables (II.1.5); then integrating over the field in each slice, an effective potential for the sum of slices lower than a given one is obtained and is renormalized in an inductive way. This is a beautiful formalism, the most natural from the point of view of the multiscale representation, and it avoids some of “the combinatoric mess to relate the counterterms to the individual graphs” [Ros]. The main technical combinatoric tool, the Gallavotti-Nicoló tree, plays however the same organizing rôle as Zimmermann forests or the tree structure of almost local subgraphs; hence everybody may agree that the key technical aspect of renormalization lies in such a structure (recall that a rooted tree minus its root is a forest, so that the difference is extremely tiny; in particular it is simply a mistake induced by their names to believe that forests are more complicated than trees...).

Here we refer the reader to the above litterature but choose to still explain the basis of the Bogoliubov recursion and of the forest formula in the old fashioned way. We do not want however to focus too much attention on them, first because it is by now rather standard material, but also because in the broader context of constructive theory, the renormalized series are not the relevant objects anyway; they are rather a natural dead end to explore before to pass to the effective series and their constructive generalizations.

The Bogoliubov induction is on the size of the graphs; it defines the counterterms to be associated to each connected superficially divergent graph at order  $n+1$  when the same thing has been done up to order  $n$ . Most of the “combinatoric mess” is simply due to the use of symmetry factors  $S(G)$  in formulas for graphs and disappears if one returns instead to the notion of contraction schemes, which we abbreviate in this section as CS. Also the presence of an ultraviolet cutoff is necessary for well defined formulas and we assume it.

The best thing is before the general rule to gain some practice with simple examples derived from the perturbation theory of the 4 point function, in which

we do not consider graphs which contain subgraphs with  $N = 2$  (called “bipeds” in [dCR1]). The lowest-order divergent graph is then the bubble. It corresponds to many CS, several for each of the three traditional  $s$ ,  $t$  and  $u$  channels of scattering theory. We introduce a counterterm in the Lagrangian  $c_{\text{bubble}}\phi^4$  with  $c_{\text{bubble}}$  defined by the subtraction prescription; in the BPHZ scheme it is minus the value of the bare amplitude at 0 momenta, but the principle of the Bogoliubov recursion would be the same for other subtraction schemes. With this new term in the Lagrangian, the four point function becomes finite up to second order. What is less obvious is that for each graph containing a single bubble subgraph, there is exactly one associated subgraph containing a single vertex of type  $c_{\text{bubble}}$ ; the combinatoric to check, at the level of CS, is simply the multinomial formula which allows to choose a pair of vertices among  $n$  and to build the bubble with them. Similarly for a graph with several disjoint bubbles one new contribution is generated for each non empty subset of these bubbles, where each bubble in the subset is replaced by  $c_{\text{bubble}}$ . Again the combinatoric is checked through the multinomial formula for the (unordered) choice of several pairs of vertices among  $n$ . For instance at third order we get the contributions of Fig.II.2.5.

Grouping together the contributions in the natural way we obtain partially renormalized amplitudes  $A^{PR}$ . They are pictured in a symbolic, but hopefully transparent way in Fig.II.2.6.

These partially renormalized amplitudes are not convergent when the cutoff is removed, but partial subintegrations previously divergent are now convergent. We introduce a new third order counterterm  $c_G$  for each third order graph  $G$  which is minus the value of  $A_G^{PR}$  at 0 external momenta. With this new Lagrangian, the amplitudes associated to the graphs of Fig.II.2.5-6 become the renormalized ones  $A_G^R = A_G^{PR} + c_G$ . Again after some use of multinomial coefficients new structures appear at higher order which contain insertions of third order counterterms for each possible set of reductions of a third order divergent subgraph to a single vertex. This notion will be made soon more precise by the use of forests, but for the moment we invite the reader to try still particular examples and to check the multinomial coefficients.

We can now state the general principle of the Bogoliubov induction:

$$c_G = - \sum_{g_1, \dots, g_k} \tau_G A_{G/\{g\}} \prod_{i=1}^k c_{g_i}|_{p=0} \quad (\text{II.2.11})$$

which means:

- For each possible family  $\{g\} = \{g_1, \dots, g_k\}$  of *disjoint* connected superficially divergent subgraphs  $g_i$  of  $G$ , multiply the corresponding counterterms, which by induction have been defined at an earlier stage,
- multiply for each such family by the bare amplitude for the “reduced graph”  $G/\{g\}$  obtained by reducing each  $g_i$  to a single vertex (of the correct type, i.e. with  $N(g_i)$  lines hooked to it),
- take minus the beginning of a Taylor expansion in the external momenta (at the value 0 in this subtraction scheme) as indicated by the operator  $\tau_G$  and sum over all possible such families, including the empty one.

We can remark that we may use other subtraction prescriptions, for we may oversubtract; for instance in the  $\phi_3^4$  theory we could use the same scheme as for

$\phi_4^4$  and introduce counterterms for the coupling constant, although they would be in this case finite as the cutoff is removed, and therefore not necessary from the point of view of power counting. In this way the Bogoliubov recursion treats finite or infinite renormalizations on the same footing.

The renormalized functional integral is then formally given by:

$$\frac{1}{Z} e^{-(1/4!)(g_r + \sum_n g_r^n c_n) \int \phi^4 - (1/2)(m^2 + \sum_n g_r^n d_n) \int \phi^2 - (1/2)(a + \sum_n g_r^n e_n) \int \partial_\mu \phi \partial^\mu \phi} D\phi \quad (\text{II.2.12})$$

where  $c_n$  is given by the sum over graphs  $G$  with  $n$  vertices and 4 external legs (the “quadrupeds”), of the counterterm  $c_G$  of (II.2.11); and  $d_n$  and  $e_n$  correspond respectively to the first and the second subtraction of the Taylor operator  $\tau_G$  for graphs  $G$  with  $n$  vertices and two external legs (the “bipeds”).

We insist on the purely formal character of the functional integral (II.1.14). It is worse than the bare functional integral (I.3.1), which becomes the well defined functional integrals (I.3.9) or (I.3.12) when cutoffs are applied. Even with cutoffs, (II.1.14) remains ill defined because the infinite series  $c_n$ ,  $d_n$  and  $e_n$  are formal and not expected to be convergent.

If one accepts this enormous drawback, we are now in a position to check Zimmermann’s forest formula, which derives the renormalized perturbation series from (II.1.14). The result is given in terms of renormalized amplitudes in which a certain set of subtractions is performed directly on the *integrand*, so that the resulting integrals are absolutely convergent. The advantage is that this derivation, although formal, works as well for the theory with and without ultraviolet cutoff.

The renormalized amplitudes  $A_G^R$  in [Zim] are then similar to the bare amplitudes in momentum space (I.4.18) but with the insertion of a renormalization operator:

$$A_G^R(p_1, \dots, p_N) = \int \mathbf{R} \prod_{l \in G} \frac{d^4 p_l}{p_l^2 + m^2} \prod_{v \in G} \delta\left(\sum_l \epsilon_{v,l} p_l\right) \quad (\text{II.2.13})$$

$$\mathbf{R} = \sum_{\mathbf{F}} \prod_{g \in \mathbf{F}} (-\tau_g) \quad (\text{II.2.14})$$

The sum runs over all possible forests  $\mathbf{F}$  of subgraphs which are connected and superficially divergent, hence all possible forests of quadrupeds and bipeds, *including the empty one*, which corresponds to the bare amplitude. We recall that a forest is a set of subgraphs so that any two elements are either disjoint or included one in the other. The reader is invited to look for some examples, for instance to the twelve forests which contribute to the sum for the graph  $G$  in Fig.II.2.7.

The BPHZ scheme (II.2.13-14) is characterized by the following normalization conditions on the connected functions in momentum space:

$$C^4(0, 0, 0, 0) = -g_r \quad (\text{II.2.15})$$

$$C^2(p^2 = 0) = \frac{1}{m^2} \quad (\text{II.2.16})$$

$$\frac{d}{dp^2} C^2|_{p^2=0} = -\frac{a}{m^4} \quad (\text{II.2.17})$$

These conditions are often stated in terms of the vertex functions (one particle irreducible amputated functions).

We sketch now the derivation of (II.2.13-14) from (II.2.12). The fact that forests are the solution to Bogoliubov's recursion is quite obvious, because considering families of disjoint subgraphs among the members of a family of disjoint subgraphs, etc... produces obviously forests. We have to check that each forest is produced exactly once for each graph, at the level of Wick contractions. The corresponding combinatoric to check reduces to the multinomial formula for choosing the sequence of  $n_i$  vertices  $i = 1, \dots, k$  among  $n$  in each maximal subgraph  $g_i$  in the forest (the “root” of each tree, from the point of view of the inclusion relation). The combinatoric inside each maximal such subgraph has not to be checked since it is treated automatically by the induction.

Therefore the last subtle point is to check that in formula (II.2.13) the **R** operator may be taken to act really on the integrand. This was accomplished in [Zim]. In (II.2.13-14) the elementary operators  $\tau_g$  for instance for forests made of a single subgraph  $g$  are really given by (II.2.2), namely for logarithmically divergent subgraphs the operator  $\tau_g$  simply puts to zero their external momenta; for more divergent subgraphs a Taylor expansion around zero momenta is taken as in (II.2.2). However when we want to combine these elementary operations into products over subgraphs of a forest some technical ambiguities have to be fixed, and this makes the true definition of the **R** operator in (II.2.13-14) rather complicated. The product of the Taylor operators (II.2.2) has to be applied in the natural order of the forest, starting from the smallest graphs (the “leaves”). But one has also in fact to eliminate the  $\delta$  functions in (II.2.13) by choosing a momentum routing rule which must be “admissible”; then one can define the action of the **R** operator on the reduced integrand and check that the result is independent of the admissible routing chosen. The insertion of the **R** operator leads in the end to well defined absolutely convergent integrals [Zim].

These complications were bypassed in the parametric representation, where Bergère, Lam and Zuber [BL],[BZ] proved that there is an equivalent but completely canonical definition of the **R** operator, which also acts by direct subtractions on the integrand, and makes the Feynman integrals absolutely convergent, provided one simply works in the  $\alpha$  representation (I.4.19) rather than in the momentum representation (I.4.18). Moreover the Taylor operators defined in  $\alpha$  space may be shown to commute precisely when the graphs belong to a common forest, so that there is no ordering ambiguity in formula (II.1.18) in this case. This formalism, which was used in [dCR1][Ri1], is certainly the most elegant and compact one for a purely perturbative definition of renormalization.

However we want neither to work in the momentum, nor in the parametric representation, but in phase space. Hence we will define and use below a third equivalent version in which the **R** operator acts in  $x$ -space. In this formalism, which was developed in [FMRS2], the momentum space Taylor operations  $\tau_G$  must be replaced by  $x$ -space “adjoints”  $\tau_G^*$ , which generalize (II.2.3-4) and whose exact definition is given in the next section. This definition, like (II.2.3-4), implies the non-canonical choice of one particular vertex  $v_e$  among the border vertices of each subgraph of a given forest; there is also a condition that this non-canonical choice must be done coherently throughout the forest, which can be considered a

dual of the admissibility condition on Zimmermann's routings of momenta. This sounds like a backward-step with respect to the  $\alpha$  space formalism of [BL][BZ]; but the advantage is that this formalism is fully compatible with our multiscale representation; therefore it can be extended naturally to constructive field theory.

Starting from such a well defined formula for renormalized amplitudes, our goal is to show not only that these amplitudes are indeed finite (the “BPH theorem”) but again to find also good uniform bounds at large order. The simplest such bound is obtained for amplitudes with fixed external momenta, but of course it is possible to derive analog results for amplitudes with external points smoothed against given test functions (see H1-H3 in the preceding section). Let us define  $f(G)$  as the supremum over all forests appearing in (II.2.14) of  $|\mathbf{F}|$ , the number of subgraphs in  $\mathbf{F}$ . It is easy to check that  $f(G) \leq n(G)$ . As remarked already, a bound similar to (II.1.9) is hopeless because of the renormalon factorials. Instead we want to prove:

### **Theorem II.2.1: The BPH uniform theorem**

There exists some constant  $K$  such that:

$$A_G^R(p_1, \dots, p_N) \leq K^n f(G)! (1 + \sup_j |p_j|)^{\hat{N}} \quad (\text{II.2.18})$$

where  $\hat{N}$  is a function of  $N(G)$  which may be taken to be  $\hat{N} = N(G)/2$  if  $N \geq 6$ ,  $\hat{N}(G) = 1$  if  $N(G) = 4$  and  $\hat{N}(G) = 3$  if  $N(G) = 2$ .

This bound is more accurate than those of [dCR1] and [FMRS2] as far as the external momenta dependence is concerned. (II.2.18) should not be considered to reflect the true large momentum behavior of  $I_G^R$ . For quadrupeds for instance one may establish as an exercise a more precise bound like

$$\sum_{k=1}^{f(G)} (f(G) - k)! [\log(1 + \sup_i |p_i|)]^k \quad (\text{II.2.19})$$

for the right hand side of (II.2.18).

Nevertheless, the bound (II.2.18) for the first time gives a radius of convergence in the Borel plane which is uniform in the external momenta. This is an improvement on [FMRS2] due to our use of the replica trick (II.1.22). Even without this trick it should be possible to get a Borel radius uniform in the external momenta, but the analysis is more tedious and we do not try it here.

When later combined with an analysis of the number of graphs with fixed  $n$  and  $f$ , (II.2.18) leads to a finite “Borel radius” of the renormalized series, and one which does not shrink at large momenta like those of [dCR1], [FMRS2]. An improved analysis sketched in Sect.II.6 even leads to bounds which give the optimal expected radius of convergence in the Borel plane. For the moment however we do not try to find optimal constants  $K$  in (II.2.18).

The next section is devoted to a full proof of the theorem. For pedagogical reasons it is divided in two steps; the proof is given first for graphs without bipeds (connected 2 point subgraphs), then in the general case. Indeed bipeds of  $\phi_4^4$  are the source of technical complications. To avoid redundant subtractions which result in worse bounds, it is natural to push for them the analysis beyond the level of connected subgraphs, to the level of one-particle irreducible subgraphs.

The natural forests associated with this notion are the forests of “closed graphs” [dCR1]. However the corresponding technicalities divert the attention from the core of the proof, which appears more clearly in the biped-free case. Hence the hope is that with this presentation, the reader will *not* skip at least the first part of the next section.

### II.3 Proof of the uniform BPH theorem

#### A. The biped-free case

For the proof to work in the multiscale representation, the key point is to extract additional “index space” decay when some almost local subgraphs are divergent; in the biped free case they can only be quadrupeds. This decay should be extracted as in the example of the bubble subgraph treated in the preceding section. We need to give first the precise definition of the subtraction operators which are the equivalents in  $x$ -space of Zimmermann’s operators in momentum space; in other words we must give the rule for “adjoints”  $\tau_g^*$  which generalize (II.2.3-4). Let us fix some forest of quadrupeds  $\mathbf{F}$ . We want to define the equivalent of the product  $\prod_{g \in \mathbf{F}} \tau_g$  acting on the momentum space integrand  $\prod_l (p_l^2 + m^2)^{-1} \prod_v \delta(\sum \epsilon_v, l p_l)$  as a product  $\prod_{g \in \mathbf{F}} \tau_g^*(v_e(g))$ , acting on the  $x$ -space integrand, so that the renormalized amplitudes computed with the  $x$ -space  $\tau^*$  operators agree with the ones computed with Zimmermann’s  $\tau$  operators. (Since  $C(x, y)$  is ill defined at coinciding points, we define in fact the  $\tau^*$  operators as acting always on a regularized integrand  $\prod_l C_\kappa(x_l, y_l)$ , and the preceding statement will be true only in the limit  $\kappa \rightarrow \infty$ .) The definition of the  $\tau^*$  operators simplifies in our case because we meet only quadrupeds, for which  $\omega = 0$ . Therefore the sum (II.2.2) reduces to a single term, the 0 momentum value. We show now that in  $x$ -space, the equivalent of taking external momenta to 0 is to integrate over the position of vertices, save one (which corresponds to global translation invariance, hence to the overall  $\delta$  function of momentum conservation for the subgraph in Zimmermann’s scheme). Hence we must do a consistent choice for all the subgraphs  $g$  of a given forest of quadrupeds  $\mathbf{F}$  of a preferred or “fixed” border-vertex  $v_e(g, \mathbf{F})$  and define the corresponding  $\tau_g^*(v_e(g, \mathbf{F}))$  operators. The following rule is a correct one (not unique).

Choose a border vertex arbitrarily for any of the maximal subgraphs of  $\mathbf{F}$  (the trunks), but one which, if possible, is also a border vertex for  $G$  itself. Then choose inductively the other border vertices according to the natural rule: if  $g'$  is the immediate ancestor of  $g$  in the forest, which we note  $g' = B_{\mathbf{F}}(g)$  and  $v_e(g', \mathbf{F})$  is also a border vertex of  $g$ , choose  $v_e(g, \mathbf{F}) = v_e(g', \mathbf{F})$ . If it is not the case but there are some border vertices of  $g$  which are also border vertices of  $g'$ , choose  $v_e(g, \mathbf{F})$  among them; otherwise choose  $v_e(g, \mathbf{F})$  arbitrarily.

With this simple rule we may picture in a graphic way our definition of the action of the product  $\prod_{g \in \mathbf{F}} \tau_g^*(v_e(g, \mathbf{F}))$  on the integrand  $\prod_l C_\kappa(x, y) \prod_{j=1}^N e^{i \sum p_j x_j}$ . Each  $\tau_g^*$  operation simply changes every external line  $C(x, z)$  of  $g$  into  $C(x_{v_e(g)}, z)$  which means that it moves each external line of  $g$  to attach it to the single border vertex  $v_e(g, \mathbf{F})$ . These operations are consistent and commute because whenever  $g \subseteq g'$  our rule ensures that an external line common to  $g$  and  $g'$  is never moved by the  $\tau_g^*$  operator to an inside vertex of  $g'$  (which would be bad because in the definition of the  $\tau'_g$  operator the inside vertices of  $g'$  have to be integrated out and one could not apply both the  $\tau_g^*$  and  $\tau_{g'^*}$  operators). With our rule the product  $\prod_{g \in \mathbf{F}} \tau_g^*(v_e(g, \mathbf{F}))$  results in a well defined set of “moves” for the lines of  $g$ . Taken literally as in Fig.II.3.1 these moves would lead to unpleasant vertices

with more than 4 lines attached to them so it is better to view the  $\tau^*$  operators as creating thick “reduction” vertices which are then expanded separately to show their “inside” (the subgraph to which they correspond) as in Fig.II.3.2 where arrows show which “inside” is associated with which vertex.

Returning to phase space, we apply the decomposition (II.1.8) into momentum assignments and obtain:

$$A_G^R(p_1, \dots, p_N) = \sum_{\mu} \int \prod_v dx_v \mathbf{R} \prod_l C^{i_l(\mu)}(x_l, y_l) \prod_{j=1}^N e^{ip_j x_j} \quad (\text{II.3.1})$$

$$\mathbf{R} \equiv \sum_{\mathbf{F}} \prod_{g \in \mathbf{F}} [-\tau_g^*(v_e(g, \mathbf{F}))] \quad (\text{II.3.2})$$

(for simplicity we may often forget to write the dependence in  $v_e(g, \mathbf{F})$ ).

Suppose now that we have a fixed assignment  $\mu$ . The basic problem is to organize the forests appearing in the definitions (II.2.17) or (II.3.2) of  $\mathbf{R}$  according to the fundamental tree structure induced by  $\mu$ , namely the tree of the almost local  $G_k^i$ 's. The almost local divergent subgraphs form a particularly obvious subforest of this tree, which we call  $\mathbf{D}_\mu$ , the “dangerous” forest for  $\mu$  (since we remarked already at length that it is responsible for all the ultraviolet divergences in the assignment  $\mu$ ). Recalling (II.1.12), this forest may be characterized as the set of quadrupeds  $g$  with  $i_g(\mu) > e_g(\mu)$  (see (II.1.10-11) for the definition of these indices). In (II.3.2) the sum over all forests which are subforests of  $\mathbf{D}_\mu$ , reconstructs exactly the operator  $\prod_{g \in \mathbf{D}_\mu} (1 - \tau_g^*)$ . Intuitively we may guess that this is exactly what we need to restore vertical (index space) decay and to cure the ultraviolet divergences. But why are there other forests in (II.3.2), and what should we do with them?

We saw that graphically the  $\tau_g^*$  operators extract the subgraph  $g$  from  $G$  and replace it by a reduction vertex. Hence intuitively these operators should prevent any contact across the boundary of  $g$ . Therefore when an operator  $\tau_g^*$ , or more precisely a forest of such operators is applied we should try to define a new version of  $\mathbf{D}_\mu$  relative to it, which respects this constraint of no information passing through any boundary of any  $g$ . This natural idea is formalized as follows.

Given any quadruped forest  $\mathbf{F}$  and any  $g$  compatible with it (such that  $\mathbf{F} \cup g$  is still a forest) we define  $B_{\mathbf{F}}(g)$  as the ancestor of  $g$  in  $\mathbf{F} \cup g$  and  $A_{\mathbf{F}}(g)$  as  $\cup_{h: g \supset h \in \mathbf{F}} h$ .

Then for  $\mu$  and  $\mathbf{F}$  given, we define two subforests of  $\mathbf{F}$ , called respectively the safe and dangerous parts of  $\mathbf{F}$ . The safe part  $\mathbf{S}_\mu(\mathbf{F})$  of  $\mathbf{F}$  is the complement in  $\mathbf{F}$  of the dangerous part  $\mathbf{D}_\mu(\mathbf{F})$ , which is defined by:

$$g \in \mathbf{D}_\mu(\mathbf{F}) \Leftrightarrow g \in \mathbf{F}, i_g(\mathbf{F}) > e_g(\mathbf{F}) \quad (\text{II.3.3})$$

where

$$i_g(\mathbf{F}) = \min\{i_l(\mu) | l \in g - A_{\mathbf{F}}(g)\} \quad (\text{II.3.4})$$

$$e_g(\mathbf{F}) = \max\{i_l(\mu) | l \in E(g) \cap B_{\mathbf{F}}(g)\} \quad (\text{II.3.5})$$

$E(g)$  being the set of external lines of  $g$  internal in  $G$ .  $E(g) \cap B_{\mathbf{F}}(g)$  is empty only when  $g = G$ , in which case we set  $e_g(\mathbf{F}) = -1$ . Remark that these definitions indeed generalize (II.1.10-11) since  $i_g = i_g(\emptyset)$  and  $e_g = e_g(\emptyset)$ . Remark also that

the dangerous graphs in  $\mathbf{D}_\mu(\mathbf{F})$  are exactly those  $g$  in  $\mathbf{F}$  such that  $g/A_{\mathbf{F}}(g)$  is a connected component of  $[B_{\mathbf{F}}(g)/A_{\mathbf{F}}(g)]^i(\mu)$  for some  $i$ , hence they generalize the notion of almost locality to the case of the reduced graph  $B_{\mathbf{F}}(g)/A_{\mathbf{F}}(g)$ . We recall that the notation  $g/h$  means that in  $g$  every connected component of  $h$  has been reduced to a single vertex ,and that  $g^i(\mu) = \{l \in g | i_l(\mu) \geq i\}$ .

Obviously if  $\mathbf{F}_1 \subseteq \mathbf{F}_2$  we have:

$$i_g(\mathbf{F}_1) \leq i_g(\mathbf{F}_2) ; e_g(\mathbf{F}_1) \geq e_g(\mathbf{F}_2) \quad \forall g \in \mathbf{F}_1 \quad (\text{II.3.6})$$

We have also for any  $g \in \mathbf{F}$ :

**Lemma II.3.1**

$$i_g(\mathbf{F}) = i_g(\mathbf{S}_\mu(\mathbf{F}) \cup \{g\}) \quad (\text{II.3.7})$$

$$e_g(\mathbf{F}) = e_g(\mathbf{S}_\mu(\mathbf{F}) \cup \{g\}) \quad (\text{II.3.8})$$

**Proof [FMRS2]**

Suppose  $A_{\mathbf{S}_\mu(\mathbf{F})}(g) \subset A_{\mathbf{F}}(g)$  and let  $l_0$  be any line in  $A_{\mathbf{F}}(g)/A_{\mathbf{S}_\mu(\mathbf{F})}(g)$ . Let  $d_0 \subset d_1 \dots \subset d_n$  be the set of all dangerous elements of  $\mathbf{F}$  containing  $l_0$  and contained in  $g \equiv d_{n+1}$ . For each  $k = 0, \dots, n$  there must exist a line  $l_{k+1}$  which is an external line of  $d_k$  and is contained in  $d_{k+1}$ . Since each  $d_k$  is dangerous,  $i_{l_k} > i_{l_{k+1}}$ , hence  $i_{l_0} > i_{l_{n+1}}$ , with  $l_{n+1} \in g/A_{\mathbf{F}}$ . Hence  $l_0$  cannot bear the minimal index in  $g/A_{\mathbf{S}_\mu(\mathbf{F})}(g)$ , which proves (II.3.7).

Similarly if  $E(g) \cap B_{\mathbf{F}}(g) \subset E(g) \cap B_{\mathbf{S}_\mu(\mathbf{F})}(g)$ , let  $l_0$  be any line in the difference, and  $l_1$  be any line in  $E(g) \cap B_{\mathbf{F}}(g)$ . Let  $d$  be the largest element of  $\mathbf{F}$  such that  $l_0 \in E(d)$  and  $g \subset d \subset B_{\mathbf{S}_\mu(\mathbf{F})}(g)$ .  $d$  must be dangerous for  $\mathbf{F}$  so using (II.3.8):

$$i_{l_0} \leq e_d(\mathbf{F}) < i_d(\mathbf{F}) = i_d(\mathbf{S}_\mu(\mathbf{F}) \cup \{d\}) \leq i_{l_1} \quad (\text{II.3.9})$$

hence  $l_0$  cannot provide the maximum in  $E(g) \cap B_{\mathbf{S}_\mu(\mathbf{F})}(g)$ .

As a consequence of the Lemma,  $\mathbf{S}_\mu(\mathbf{S}_\mu(\mathbf{F})) = \mathbf{S}_\mu(\mathbf{F})$  and the set  $F^D(G)$  of all quadruped forests decomposes according to classes under the action of the  $\mathbf{S}_\mu$  projector:

$$F^D(G) = \bigcup_{\mathbf{F} \mid \mathbf{S}_\mu(\mathbf{F}) = \mathbf{F}} \{\mathbf{F}' \mid \mathbf{S}_\mu(\mathbf{F}') = \mathbf{F}\} \quad (\text{II.3.10})$$

The forests  $\mathbf{F}$  satisfying  $\mathbf{S}_\mu(\mathbf{F}) = \mathbf{F}$  form the set  $\mathbf{Safe}(\mu)$  of the so called safe forests (with respect to  $\mu$ ).

For any such safe forest  $\mathbf{F}$ , the next lemma characterizes completely the equivalence class  $\{\mathbf{F}' \in F^D(G) \mid \mathbf{S}_\mu(\mathbf{F}') = \mathbf{F}\}$  in terms of the forest:

$$\mathbf{H}_\mu(\mathbf{F}) = \{g \subseteq G \mid g \text{ quadruped compatible with } \mathbf{F} \text{ and } g \in \mathbf{D}_\mu(\mathbf{F} \cup \{g\})\} \quad (\text{II.3.11})$$

**Lemma II.3.2** For any  $\mathbf{F} \in \mathbf{Safe}(\mu)$  one has:

$$\mathbf{F} \cup \mathbf{H}_\mu(\mathbf{F}) \in F^D(G) \quad (\text{II.3.12a})$$

$$\forall \mathbf{F}' \in F^D(G), \mathbf{S}_\mu(\mathbf{F}') = \mathbf{F} \Leftrightarrow \mathbf{F} \subseteq \mathbf{F}' \subseteq \mathbf{F} \cup \mathbf{H}_\mu(\mathbf{F}) \quad (\text{II.3.12b})$$

**Proof [FMRS2]**

(a) We have to show that two subgraphs  $g, g'$  in  $\mathbf{H}_\mu(\mathbf{F})$  cannot overlap, i.e. have a nontrivial intersection. If they do, since  $g$  is connected,  $g - g'$  must contain a line  $l$  which is internal for  $g$  and external for  $g'$ ; this line  $l$  cannot be in  $A_{\mathbf{F}}(g)$  (otherwise  $g'$  would not be compatible with  $\mathbf{F}$ ) and must be in  $B_{\mathbf{F}}(g')$  (by compatibility of  $g$  with  $\mathbf{F}$ ). Similarly  $g'$  must contain a line  $l'$  in  $g'/A_{\mathbf{F}}(g') \cap E(g) \cap B_{\mathbf{F}}(g)$ . This enforces a contradiction:

$$i_l(\mu) \geq i_g(\mathbf{F} \cup \{g\}) > e_g(\mathbf{F} \cup \{g\}) \geq i_{l'}(\mu); i_{l'}(\mu) \geq i_{g'}(\mathbf{F} \cup \{g'\}) > e_{g'}(\mathbf{F} \cup \{g'\}) \geq i_l(\mu) \quad (\text{II.3.13})$$

(b)  $\Leftarrow$  Note first that if  $\mathbf{F}_1, \mathbf{F}_2 \in F^D(G)$  and  $\mathbf{F}_1 \subseteq \mathbf{F}_2$ , then by (II.3.6):

$$g \in \mathbf{D}_\mu(\mathbf{F}_1) \Rightarrow g \in \mathbf{D}_\mu(\mathbf{F}_2) \quad (\text{II.3.14})$$

Hence if  $g \in \mathbf{F}' - \mathbf{F} = \mathbf{F}' \cap \mathbf{H}_\mu(\mathbf{F})$ , then  $g \in \mathbf{D}_\mu(\mathbf{F} \cup \{g\}) \subseteq \mathbf{D}_\mu(\mathbf{F}')$ . So  $\mathbf{S}_\mu(\mathbf{F}') \subseteq \mathbf{F}$ . If this inclusion was a strict one, it would mean that there is a  $g \in \mathbf{F}$  with  $g \in \mathbf{D}_\mu(\mathbf{F}')$ , hence using Lemma II.3.1:

$$i_g(\mathbf{F}') > e_g(\mathbf{F}') = e_g(\mathbf{S}_\mu(\mathbf{F}' \cup g)) \quad (\text{II.3.15})$$

and using (II.3.6) and the fact that  $\mathbf{F} = \mathbf{S}_\mu(\mathbf{F})$ :

$$e_g(\mathbf{S}_\mu(\mathbf{F}' \cup g)) \geq e_g(\mathbf{F}) \geq i_g(\mathbf{F}) = i_{l_0}(\mu) \quad (\text{II.3.16})$$

Combining (II.3.15) and (II.3.16), we see that the line  $l_0$  must have collapsed under reduction of  $A_{\mathbf{F}'}(g)$  but not under reduction of  $A_{\mathbf{F}}(g)$ . Therefore  $l_0$  must belong to some  $h$  in  $\mathbf{H}_\mu(\mathbf{F})$  with  $h \subset g$ , but cannot belong to  $A_{\mathbf{F}}(g)$ . Since  $h \cap A_{\mathbf{F}}(g) = A_{\mathbf{F}}(h)$  we have:

$$i_{l_0}(\mu) = \min\{i_l | l \in h/A_{\mathbf{F}}(h)\} > \max\{i_l | l \in E(h) \cap B_{\mathbf{F}}(h)\}. \quad (\text{II.3.17})$$

This contradicts the last part of (II.3.16) because  $E(h)$  must contain a line that is internal to  $B_{\mathbf{F} \cup \{h\}}(h) \subseteq g$  (or else  $B_{\mathbf{F} \cup \{h\}}(h)$  would not be connected) and this line cannot be in  $A_{\mathbf{F}}(g)$ , since  $h$  is compatible with  $\mathbf{F}$  and is not contained in  $A_{\mathbf{F}}(g)$ .

$\Rightarrow$  If  $\mathbf{S}_\mu(\mathbf{F}') = \mathbf{F}$  then  $\mathbf{F} \subseteq \mathbf{F}'$  and using Lemma II.3.1:

$$g \in \mathbf{F}' - \mathbf{F} \Rightarrow g \in \mathbf{D}_\mu(\mathbf{F}') \Rightarrow g \in \mathbf{D}_\mu(\mathbf{S}_\mu(\mathbf{F}') \cup \{g\}) \Rightarrow g \in \mathbf{H}_\mu(\mathbf{F}) \quad (\text{II.3.18})$$

Lemma II.3.2 allows to reorganize the  $\mathbf{R}$  operator in the assignment  $\mu$  as:

$$\mathbf{R} = \sum_{\mathbf{F} \in \mathbf{Safe}(\mu)} \prod_{g \in \mathbf{F}} (-\tau_g^*) \prod_{h \in \mathbf{H}_\mu(\mathbf{F})} (1 - \tau_h^*) \quad (\text{II.3.19})$$

This is a beautiful rearrangement (different for each  $\mu$ ) because the product  $\prod_{g \in \mathbf{H}_\mu(\mathbf{F})} (1 - \tau_g^*)$  will provide the desired cancellations for dangerous subgraphs within each particular counterterm of  $\mathbf{F}$  and within  $G/\mathbf{F}$  itself, and the counterterms of  $\mathbf{F}$  themselves, corresponding to safe subgraphs, will have some external scale acting as a cutoff.

To establish the uniform BPH theorem, we write (II.3.1) in the form:

$$A_G^R = \sum_{\mathbf{F} \in F^D(G)} A_{G,\mathbf{F}}^R \quad (\text{II.3.20})$$

$$A_{G,\mathbf{F}}^R \equiv \sum_{\mu | \mathbf{F} \in \mathbf{Safe}(\mu)} \int \prod_v dx_v \prod_{g \in \mathbf{F}} (-\tau_g^*) \\ \prod_{h \in \mathbf{H}_\mu(\mathbf{F})} (1 - \tau_h^*) \prod_l C^{i_l(\mu)}(x_l, y_l) \prod_{j=1}^N e^{ip_j x_j} \quad (\text{II.3.21})$$

and we give a bound of the form (II.2.18) on  $A_{G,\mathbf{F}}^R$  for every fixed  $\mathbf{F} \in F^D(G)$ . Then we can conclude because it has been shown ([dCR1], Lemma A.2) that the number of connected divergent forests in a biped-free  $\phi_4^4$  graph is at most  $8^{n(G)}$ . (Hint: bound first the number of maximal such forests. The only way quadrupeds may overlap in a biped-free graph is shown in the left of Fig.II.3.3 and by writing the 2-particle reducibility structure of  $G$  as in the right of Fig.II.3.3, one derives an inductive bound:

$$d_n \leq \sum_{p=1}^{n-1} d_p d_{n-p} \quad (\text{II.3.22})$$

for the maximal number  $d_n$  of such maximal forests over graphs  $G$  with  $n(G) = n$  (with initial condition  $d_1 = 1$ ). Hence  $d_n \leq 4^n$  and one can conclude).

We proceed to evaluate the action of the  $\tau^*$  and  $(1 - \tau^*)$  operators in (II.3.21). The  $\tau_g^*$  operators are applied first, for every  $g \in \mathbf{F}$ . They attach all external lines of  $g$  internal in  $B_{\mathbf{F}}(g)$  to a single reduction vertex whose position is  $v_e(g, \mathbf{F} \cup \mathbf{H}_\mu(\mathbf{F}))$ , and which we abbreviate as  $v_e(g)$  in what follows (see Fig II.3.1-2). Then we apply the  $(1 - \tau_h^*)$  operators for each  $h \in \mathbf{H}_\mu(\mathbf{F})$  and like in (II.2.5) we decompose the corresponding difference of products of external propagators into a sum of at most three terms, each of which contains exactly one difference concerning a single external propagator taken at two different end arguments. The generic case is:

$$\prod_{i=1}^4 C(x_i, y_i) - \prod_{i=1}^4 C(x_1, y_i) = C(x_1, y_1). \\ \cdot \sum_{i=2}^4 \prod_{2 \leq j < i} C(x_j, y_j) [C(x_i, y_i) - C(x_1, y_i)] \prod_{i < j \leq 4} C(x_1, y_j) \quad (\text{II.3.23})$$

Remark that the line bearing the difference must be in  $E(h) \cap B_{\mathbf{F}}(h)$ , because common external lines for  $h$  and  $B_{\mathbf{F}}(h)$  are all attached to the same vertex after the  $\tau_{B_{\mathbf{F}}(h)}^*$  operation, hence the corresponding difference in (II.3.23) gives 0 for this case and can be discarded.

Therefore the product of the  $(1 - \tau_h^*)$  operators factorizes as

$$\prod_{g \in \mathbf{F} \cup \{G\}} \left[ \prod_{h \in \mathbf{H}_\mu(\mathbf{F}) | B_{\mathbf{F}}(h)=g} (1 - \tau_h^*) \right],$$

each internal product acting only in  $g/A_{\mathbf{F}}(g)$ . Finally when we compute such a product, any difference  $C(x_1, y) - C(x_2, y)$  created at one end of a line by a  $(1 - \tau_h^*)$

operator cannot be modified by an other  $(1 - \tau_{h'}^*)$ . This is because with our rule for the choice of external vertices, the  $\tau_{h'}^*$  operator would either let both  $C(x_1, y)$  and  $C(x_2, y)$  untouched or move them both to  $C(x_{v_e(h')}, y)$ , and the difference would give 0 and can be discarded. Of course this concerns only one given end of a line, and at the other end an other difference may appear (this would be the case for Fig.II.3.1-2). Hence after all operators  $\tau_g^*$  and  $(1 - \tau_h^*)$  have been applied we get an integrand similar to the initial one, with ordinary lines (some ends of which may have been moved to some  $v_e(g)$  vertex), and lines bearing one difference at one end or two differences, one at both ends. The lines with one difference are then replaced by the interpolating formula (II.2.6) and the lines with two differences by a similar formula:

$$\begin{aligned} & C^i(x, y) - C^i(u, y) - C^i(x, z) + C^i(u, z) \\ &= \int_0^1 \int_0^1 dt_1 dt_2 \frac{d}{dt_1} \frac{d}{dt_2} C^i(x + t_1(u - x), y + t_2(y - z)) \end{aligned} \quad (\text{II.3.24})$$

Some of the lines bearing a (single!) difference may be true external lines of  $G$ ; there is one such line if  $N(G) = 4$  and at most  $N(G)/2$  if  $N(G) \geq 6$  (a pair of external lines of  $G$  could be external to some quadruped  $h$ , and one line in the pair could bear a difference; this is not true for single lines, because the rule in this case chooses the corresponding border vertex of  $h$  to be the fixed one  $v_e(h)$  and this external line does not bear any difference).

We evaluate interpolating lines of the type (II.2.6) by (II.2.7), and lines of the type (II.3.24), by a similar estimate, based on representation (II.1.3-5). A line of index  $i$  bearing 0, 1 or 2 differences will therefore have power counting  $M^{2i}$ ,  $M^{3i}$  or  $M^{4i}$  and 0, 1 or 2 corresponding multiplicative distances (like  $|x - u||y - z|$  for (II.3.24)); it will have also interpolating decay (see (II.2.7)). Before to integrate over the positions of vertices we have to bound these distance factors and replace the interpolating decays by ordinary ones, like in (II.2.8). This must be done by using the decay of the internal lines connecting  $x$  to  $u$  (or  $y$  to  $z$ ), which must be of higher scales for harvesting some net gain.

If the difference  $|x - u|$  to bound has been created by a  $(1 - \tau_h^*)$  operator with  $B_{\mathbf{F}}(h) = g$ , the line bearing the difference has index at most  $e_h(\mathbf{F})$ . The internal decay to use takes place entirely within  $g/A_{\mathbf{F}}(g)$ . But  $h$  belongs to  $\mathbf{H}_\mu(\mathbf{F})$ , which means that  $h/A_{\mathbf{F}}(h)$  is a connected component of  $(g/A_{\mathbf{F}}(g))^i$  for  $i = i_h(\mathbf{F})$ . Therefore using the triangular inequality and the decay of the internal lines of  $h/A_{\mathbf{F}}(h)$ , we can bound the difference in external arguments by  $KM^{-i_h(\mathbf{F})}$  as in (II.2.8); combining this with the extra power counting factor of the line bearing the difference results in a net gain of at least:

$$M^{-|i_h(\mathbf{F}) - e_h(\mathbf{F})|}. \quad (\text{II.3.25})$$

We can also replace the interpolating decay in  $|x + t(u - x) - y|$  by the regular one  $|x - y|$  and still keep a fraction of the internal decay of  $h/A_{\mathbf{F}}(h)$  for later use. But a subtlety arises: what if several distance factors  $|x_1 - u_1|$ ,  $|x_2 - u_2|$  etc.... made use of the decay of the same internal lines? In this case they would correspond to graphs  $h_1, h_2, \dots$ , which must be ordered by inclusion, hence have  $i_{h_1}(\mathbf{F}) > i_{h_2}(\mathbf{F}) > \dots$ . But we can duplicate the decay of any internal line by

formula (II.1.22) before using it. With this trick, different distance factors will always use different copies of any given internal decay, which solves the problem.

There is no dependence on interpolating parameters  $t$  any more and we can therefore bound the corresponding integrals by 1. This achieves the preparation of the integrand, and the rest will be very similar to Sect.II.1. We are ready to perform the integral over the positions of vertices by choosing a spanning tree  $T$  whose restriction to  $g/A_{\mathbf{F}}(g)$  is still a spanning tree  $T_g$  of  $g/A_{\mathbf{F}}(g)$  for each  $g \in \mathbf{F} \cup \{G\}$ ; this is possible because of the tree structure of  $\mathbf{F} \cup \{G\}$ . Moreover we require that for each  $i, k$  the restriction of  $T_g$  to  $(g/A_{\mathbf{F}}(g))_k^i$  is again a spanning tree of  $(g/A_{\mathbf{F}}(g))_k^i$ . This is the natural generalization of the choice of  $T$  in Sect.II.1, and is possible because of the tree structure of the  $(g/A_{\mathbf{F}}(g))_k^i$  for each  $g$ . Then the spatial integration inside each  $g/A_{\mathbf{F}}(g)$  goes exactly as for the full  $G$  in Sect.II.1, and delivers altogether the factor

$$\prod_{g \in \mathbf{F} \cup \{G\}} \prod_{(i,k)} M^{-\omega(g/A_{\mathbf{F}}(g))_k^i} \quad (\text{II.3.26})$$

Moreover the combined effect of the extra factors (II.3.25) adds to this estimate a factor which is at worst:

$$\prod_{h \in \mathbf{H}_\mu(\mathbf{F})} M^{-|i_h(\mathbf{F}) - e_h(\mathbf{F})|}. \quad (\text{II.3.27})$$

More precisely we must take into account the case where external lines of  $G$  bear differences, where we write  $|e^{ip_j x} - e^{ip_j u}| \leq (1 + \sup_j |p_j|)|x - u|$  instead of the regular estimate.  $|x - u|$  is a distance factor similar to the previous one, and the bound after integration of internal vertices is therefore:

$$|A_{G,\mathbf{F}}^R(p_1, \dots, p_N)| \leq (1 + \sup_j |p_j|)^{\hat{N}} \sum_{\mu | \mathbf{F} \in \mathbf{Safe}(\mu)} \prod_{g \in \mathbf{F} \cup \{G\}} \prod_{(i,k)} M^{-\omega'((g/A_{\mathbf{F}}(g))_k^i)} \quad (\text{II.3.28})$$

where

$$\omega'((g/A_{\mathbf{F}}(g))_k^i) \equiv \sup\{1, \omega((g/A_{\mathbf{F}}(g))_k^i)\} \quad (\text{II.3.29})$$

except when  $g \in \mathbf{F}$  and  $(g/A_{\mathbf{F}}(g))_k^i = g/A_{\mathbf{F}}(g)$ , in which case:

$$\omega'((g/A_{\mathbf{F}}(g))_k^i) \equiv \omega((g/A_{\mathbf{F}}(g))_k^i) = 0 \quad (\text{II.3.30})$$

Let  $i_{\max}(\mu)$  be the largest index appearing in the assignment  $\mu$ . From half of the index space decay in (II.3.28) we can extract a factor  $M^{-\delta i_{\max}(\mu)}$  patching together the internal vertical decays inside each  $g/A_{\mathbf{F}}(g)$ ,  $g \in \mathbf{F} \cup \{G\}$  with the condition  $i_g(\mathbf{F}) \leq e_g(\mathbf{F})$  for each  $g \in \mathbf{F}$ . Then with the rest of the decay we can sum over every assignment of each  $g/A_{\mathbf{F}}(g)$  as in Sect II.1, provided one index of  $g/A_{\mathbf{F}}(g)$  for each  $g \in \mathbf{F}$  is kept fixed; indeed (II.3.30) implies overall vertical translation invariance for each  $g/A_{\mathbf{F}}(g)$ ,  $g \in \mathbf{F}$ . The result is bounded by:

$$\prod_{g \in \mathbf{F} \cup \{G\}} K^{n(g/A_{\mathbf{F}}(g))} \leq K^{2n(G)} \quad (\text{II.3.31})$$

Each “translation invariant” sum over the fixed index of a given  $g/A_{\mathbf{F}}(g)$  is nevertheless obviously bounded by  $i_{\max}(\mu)$ , hence the corresponding sums combined are bounded by:

$$\sum_{i_{\max}(\mu)} (i_{\max}(\mu))^{|F|} M^{-\delta i_{\max}(\mu)} \leq |F|! K^{|F|} \quad (\text{II.3.32})$$

Combining (II.3.28) (II.3.31) and (II.3.32) achieves the proof of Theorem II.2.1 in the biped free case.

In fact this proof again achieves more than the theorem it is designed for. It is indeed worth to consider as a separate theorem the particular case of an empty  $\mathbf{F}$ , because the empty forest is the only safe forest common to all assignments, and by (II.3.32) there is no factorial in the estimates for it. Since  $\mathbf{H}_\mu(\emptyset) \equiv \mathbf{D}_\mu$ , we have proved that in the biped free case:

### Theorem II.3.1: Uniform bound for usefully renormalized amplitudes

$$\begin{aligned} |A_G^{UR}(p_1, \dots, p_N)| &= \left| \sum_{\mu} \int \prod_v dx_v \prod_{h \in \mathbf{D}_{\mu}} (1 - \tau_h^*) \prod_l C^{i_l(\mu)}(x_l, y_l) \prod_{j=1}^N e^{ip_j x_j} \right| \\ &\leq (1 + \sup_j |p_j|)^{\hat{N}} K^{n(G)} \end{aligned} \quad (\text{II.3.33})$$

$A_G^{UR}$  is a piece of the renormalized amplitude, called the “usefully” renormalized amplitude because it has no insertion of “useless” counterterms which correspond to the elements of safe forests. Theorem II.3.1 shows that it does not contain any renormalon effect. Hence we can conclude, as announced in the previous section, that the renormalon effects are solely due to the useless counterterms.

We give now the proof of the BPH uniform theorem (Theorem II.2.1) in the general case where bipeds are present. We will also obtain the generalization of Theorem II.3.1, proving that it also holds with bipeds. We will not rephrase what is similar to the above analysis but concentrate on the new technicalities created by the bipeds.

## B. The general case, with bipeds

To get a reasonably accurate bound when bipeds are present, we must trim some redundant subtractions in the definition (II.2.14) of the  $\mathbf{R}$  operator due to one particle reducible divergent graphs. Indeed for a graph with only  $n$  vertices like the one of Fig.II.3.4 there are forests of connected divergent subgraphs with about  $3n/2$  elements, and the corresponding estimate (II.2.18) would not even allow a finite disk of analyticity in the Borel plane (see Sect.II.6).

We recall that a subgraph was called proper or 1PI, if it cannot be broken into two disconnected pieces by cutting a single line. A quadruped (connected subgraph  $q$  with  $N(q) = 4$ ) is called *open* if it is proper and there exists a proper biped  $b \supset q$  such that both border vertices of  $b$  are border vertices of  $q$ .  $b$  is then called the closure  $q^*$  of  $q$  and is obtained by adding to  $q$  a line or a chain of proper bipeds between these two vertices (see Fig.II.3.5). More generally a subgraph  $g$  is called closed iff:

$$\forall q \text{ open quadruped } , q \subseteq g \Rightarrow q^* \subseteq g$$

and the closure  $g^*$  of a proper subgraph  $g$  is the smallest closed subgraph containing it [dCR1]. This definition is consistent with the first one when  $g$  is a quadruped.

Then we can restrict, in the definition (II.2.14) of the  $\mathbf{R}$  operator the sum to run over forests of closed divergent subgraphs (in short “closed forests”). The fact that subtractions for one particle reducible subgraphs are redundant is rather standard (see for example [BL]); it follows intuitively from the observation that putting to 0 the external momenta of a graph also puts to 0 the external momenta of its proper parts by momentum conservation. The fact that subtractions for open quadrupeds are redundant is less well known but can be grasped as follows. Let  $q$  be an open quadruped and  $b = q^*$  be its closure. Then  $(1 - \tau_b)\tau_q = 0$  since the  $\tau_q$  operator hooks both external lines of  $b$  to the same reduction vertex, and it is therefore not necessary to introduce counterterms for both  $q$  and  $b = q^*$ . From this idea and some induction, one can restrict the sum over forests in (II.2.14) to the closed forests without actually changing the action of  $\mathbf{R}$  on a graph ([dCR1], Lemma II.3). As a general definition of the factor  $f(G)$  appearing in Theorem II.2 we take the supremum over all closed forests  $\mathbf{F}$  of  $|\mathbf{F}|$ , the number of elements in  $\mathbf{F}^*$ .

For any given closed forest  $\mathbf{F}$  and assignment  $\mu$ , we introduce relative indices which generalize (II.3.4)-(II.3.5):

$$i_g^c(\mathbf{F}) = \max_h \min\{i_l(\mu) | l \in h/A_{\mathbf{F}}(g)\} \quad (\text{II.3.34})$$

where the max is taken over all  $h$  compatible with  $\mathbf{F}$  which obey  $h^* = g$  and such that  $h/A_{\mathbf{F}}(g)$  is a proper component of  $(g/A_{\mathbf{F}}(g))^i$  for some  $i$ . The  $h$  which realizes this maximum is called  $I_g^c(\mathbf{F})$ . We define  $e_g^c$  as:

$$\begin{aligned} e_g^c(\mathbf{F}) &= \max\{i_l(\mu) | l \in E(g) \cap B_{\mathbf{S}_\mu^c(\mathbf{F})}(g)\} && \text{if } N(g) = 4 \\ e_g^c(\mathbf{F}) &= \min\{i_l(\mu) | l \in E(g) \cap B_{\mathbf{S}_\mu^c(\mathbf{F})}(g)\} && \text{if } N(g) = 2 \end{aligned} \quad (\text{II.3.35})$$

where  $\mathbf{S}_\mu^c(\mathbf{F})$  is defined inductively, starting from the largest graphs in  $\mathbf{F}$  towards the smallest ones, by the condition:

$$g \in \mathbf{S}_\mu^c(\mathbf{F}) \Leftrightarrow e_g^c(\mathbf{F}) \geq i_g^c(\mathbf{F}) \quad (\text{II.3.36})$$

$$g \in \mathbf{D}_\mu^c(\mathbf{F}) \Leftrightarrow e_g^c(\mathbf{F}) < i_g^c(\mathbf{F}) \quad (\text{II.3.37})$$

Since the definition of  $\mathbf{S}_\mu^c(\mathbf{F})$  is inductive there is no logical loophole between (II.3.35) and (II.3.36).

The lemmas necessary to classify the forests are the following generalizations of Lemmas II.3.1-2:

**Lemma II.3.3** For  $\mathbf{F}$  a closed forest of  $G$  and  $g \in \mathbf{F}$ :

$$i_g^c(\mathbf{S}_\mu^c(\mathbf{F}) \cup \{g\}) = i_g^c(\mathbf{F}) \quad (\text{II.3.38})$$

$$e_g^c(\mathbf{S}_\mu^c(\mathbf{F}) \cup \{g\}) = e_g^c(\mathbf{F}) \quad (\text{II.3.39})$$

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\* This is *not* the definition of  $f(G)$  in [dCR1], and the factorial bound of [dCR1] is therefore not optimal in this respect.

$$\mathbf{S}_\mu^c(\mathbf{S}_\mu^c(\mathbf{F})) = \mathbf{S}_\mu^c(\mathbf{F}) \quad (\text{II.3.40})$$

Furthermore for  $\mathbf{F}$  a safe forest, i.e. such that  $\mathbf{S}_\mu^c(\mathbf{F}) = \mathbf{F}$ , we define  $\mathbf{H}_\mu^c(\mathbf{F}) = \{g \subseteq G, g \text{ proper closed divergent subgraph compatible with } \mathbf{F}, \text{ and } g \in \mathbf{D}_\mu^c(\mathbf{F} \cup \{g\})\}$ . Then:

**Lemma II.3.4** Lemma II.3.2 still holds with the generalized definitions  $\mathbf{S}_\mu^c(\mathbf{F})$  and  $\mathbf{H}_\mu^c(\mathbf{F})$  replacing the former  $\mathbf{S}_\mu(\mathbf{F})$  and  $\mathbf{H}_\mu(\mathbf{F})$ .

**Sketch of proof** (for a detailed proof of these two Lemmas we refer to [dCR1]). (II.3.39) is obvious from definition (II.3.35), and with (II.3.38) it implies (II.3.40). (II.3.38) is non-trivial. By induction, it is enough to show that for any  $g' \neq g, g' \in \mathbf{D}_\mu^c(\mathbf{F})$ , we have  $i_g^c(\mathbf{F}) = i_{g'}^c(\mathbf{F} - \{g'\})$ . This is obvious except when  $g = B_\mathbf{F}(g')$ . In this last case we remark first that  $h_0 = I_g^c(\mathbf{F} - \{g'\})$  satisfies all the conditions to appear on the list over which the max is taken in (II.3.34) for  $i_g^c(\mathbf{F})$ . Since  $h_0/A_\mathbf{F}(g) \subseteq h_0/A_{\mathbf{F}-\{g'\}}$ , we have the inequality  $i_g^c(\mathbf{F}) \geq i_{g'}^c(\mathbf{F} - \{g'\})$ . The opposite inequality requires some care. An easy case is when  $I_g^c(\mathbf{F})$  is disjoint from  $g'$ ; then it appears also in the list for  $\mathbf{F} - \{g'\}$  and one concludes easily. The last possibility is  $g' \subset I_g^c(\mathbf{F})$  (since  $I_g^c(\mathbf{F})$  is compatible with  $\mathbf{F}$ ). In this case  $I_g^c(\mathbf{F})$  has to contain at least two external legs of  $g'$ , hence

$$i_g^c(\mathbf{F}) \leq e_{g'}^c(\mathbf{F}) < i_{g'}^c(\mathbf{F}) \quad (\text{II.3.41})$$

and  $h_1 = I_g^c(\mathbf{F}) - [g' - I_{g'}^c(\mathbf{F})]$  is in the list over which a maximum is taken in the definition (II.3.34) of  $i_g^c(\mathbf{F} - \{g'\})$ . Hence we conclude also that  $i_g^c(\mathbf{F}) \leq i_{g'}^c(\mathbf{F} - \{g'\})$ , which achieves the proof of (II.3.38).

Remark that Lemma II.3.3 remains also true if we modify slightly the definition of  $i_g^c(\mathbf{F})$  when  $g/A_\mathbf{F}(g)$  is the bubble graph of Fig I.1.1, defining it in this special case by:

$$i_g^c(\mathbf{F}) = \max\{i_1, i_2\} \quad (\text{II.3.42})$$

where  $i_1$  and  $i_2$  are the two indices in  $\mu$  of the two lines of  $g/A_\mathbf{F}(g)$ . This small change is not fundamental at all but is necessary to have a nice rule for the “bubble resummations” of [DFR] which will be briefly considered in Sect.II.6.\* The reason for which renormalization still works with this definition is that a single line in the bubble is enough to ensure spatial decay between the two border vertices of the bubble.

For more general graphs the real constraint on  $i_g^c$  to ensure such a decay is that there is at least a spanning tree of  $g/A_\mathbf{F}(g)$  made of lines with indices higher or equal to  $i_g^c(\mathbf{F})$ ; both (II.3.34) and (II.3.42) satisfy this rule, as does the construction in [GaN]; clearly there remains some flexibility in the choice of  $i_g^c$ , and some details are a matter of convenience. This remark applies also to the expansions of constructive theory, in which the mechanism of convergence is fundamental but some particular technical details are not.

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\* It has also some advantages over the “2nd max” definition for  $e_g^c$  that the curious reader will find in [Ri1], [DFR]; this definition was introduced also solely to allow explicit bubble resummations, but is worse than the simple rule (II.3.42).

Lemmas II.3.3-4 allow to reorganize the operator  $\mathbf{R}$  exactly as in (II.3.19). Also the number of closed divergent forests in a graph is still bounded by  $8^{n(G)}$  ([dCR1], Lemma A.2). Hence we can again restrict our attention to a single  $A_{G,\mathbf{F}}^R$  in (II.3.20).

We must make precise the choice of  $v_e(g)$  when  $g$  is a biped. Recall that in this case, renormalization is performed (i.e. a factor  $(1 - \tau_g^*)$  appears in (II.3.19)) even when one of the two external legs of the biped has an index bigger than the internal index  $i_g^c$  (see (II.3.35)). Hence we have to ensure that the  $\tau_g^*$  operator in this case applies to the external leg of the biped with lowest index  $e_g^c$ , otherwise there would be no net gain. This can be done by adding to our previous rule for choosing the  $v_e(g)$  vertices the prescription that when  $g$  is a biped the fixed vertex  $v_e(g)$  is the one to which the leg of highest index in  $\mu$  is hooked. This prescription does not interfere with the former rules because a border-vertex of a proper biped cannot be also border-vertex of any other proper closed divergent subgraph containing it, hence the choice of the fixed vertex for it was arbitrary in the former rules.

We have also to supply a new formula for the action of  $\tau_g^*$  when  $g$  is a biped; since  $\omega(g) = -2$  in this case, the Taylor expansion at 0 momenta is pushed to second order. When the corresponding adjoints are applied to the external propagators  $C(x, u)C(y, z)$  of the biped, with  $x = x_{v_e(g)}$  and  $y$  the position of the other border vertex of  $g$ , one obtains:

$$\begin{aligned} \tau_g^*[C(x, u)C(y, z)] &= C(x, u)\{C(x, z) + (y - x)^\mu \frac{\partial C}{\partial x^\mu}(x, z) \\ &\quad + \frac{1}{2}(y - x)^\mu(y - x)^\nu \frac{\partial^2 C}{\partial x^\mu \partial x^\nu}(x, z)\} \end{aligned} \quad (\text{II.3.43})$$

where partial derivatives apply to the first argument  $x$  of  $C(x, z)$ . By a parity argument the second term in the sum (II.3.43) vanishes when integrated over  $y$ , so that we may forget it and write  $\tau_g = \tau_g^0 + \tau_g^1$ .  $\tau_g^0$  is the mass counterterm, whose adjoint is defined by:

$$\tau_g^{0*}[C(y, z)C(x, u)] = C(x, z)C(x, u) \quad (\text{II.3.44})$$

and  $\tau_g^1$  is the wave function counterterm, whose adjoint is defined by:

$$\tau_g^{1*}[C(y, z)C(x, u)] = \frac{1}{2}(y - x)^\mu(y - x)^\nu \frac{\partial^2 C}{\partial x^\mu \partial x^\nu}(x, z).C(x, u) \quad (\text{II.3.45})$$

These counterterms are pictured in Fig.II.3.6. Again by parity, only the terms with  $\mu = \nu$  survive after the internal integration over  $y - x$  (holding  $x$  fixed) is performed. Finally we can remark that  $y - x$  and  $z - x$  are independent in (II.3.45), and we can use Euclidean invariance to replace each integral containing a particular  $[(y - x)^\mu]^2 \frac{\partial^2 C}{\partial x^\mu \partial x^\mu}(x, z)$  by  $(1/4)|y - x|^2 \frac{\partial^2 C}{\partial x^\mu \partial x^\mu}(x, z)$ . Writing  $\Delta = \sum_\mu \frac{\partial^2}{\partial x^\mu \partial x^\mu}$  we can therefore rewrite (II.3.45) as:

$$\tau_g^{1*}[C(y, z)C(x, u)] = \frac{1}{8}|y - x|^2 \Delta C(x, z).C(x, u) \quad (\text{II.3.46})$$

which is simpler and hopefully suggests more clearly that the wave function counterterm really corresponds to the renormalization of  $a$ , the parameter in (I.3.1)

which corresponds in the action to  $\int \partial_\mu \phi \partial^\mu \phi$  hence, integrating by parts, to  $-\int \Delta \phi \cdot \phi$ .

It is also useful to write the Taylor remainder formula:

$$(1 - \tau_g^*)[C(x, u)C(y, z)] = C(x, u) \int_0^1 \frac{(1-t)^2}{2} \frac{d^3}{dt^3} C(x + t(y-x), z) \quad (\text{II.3.47})$$

As before, our rules on the choice of fixed vertices prevents any end of line to bear the action of several  $(1 - \tau_g^*)$  operators. To the preceding argument for biped free graphs, one has to add the observation that an external line of a biped  $g \in \mathbf{F} \cup \mathbf{H}_\mu^c(\mathbf{F})$  cannot be the external line of other elements of  $\mathbf{F} \cup \mathbf{H}_\mu^c(\mathbf{F})$ , except maybe for quadrupeds  $h \subset g$ , and for these  $h$  our rule for choosing  $v_e(h)$  ensures that this external line common to  $g$  and  $h$  is hooked to  $v_e(h)$ , hence cannot bear the action of the  $(1 - \tau_h^*)$  operator.

We can evaluate the net effect of operations (II.3.43-47). In (II.3.47) we earn a factor at least  $M^{-3[i_g^c(\mathbf{F}) - e_g^c(\mathbf{F})]}$ , again using representation (II.1.3-4), and the internal decay of a spanning tree of  $g$  with lines of indices greater or equal to  $i_g^c(\mathbf{F})$ . Duplicating internal lines decay like in (II.1.22) again avoids using too many times the same decay of the same line.

Similarly evaluating (II.3.44) is neutral, and (II.3.45-46) results in a *loss* of at most  $M^{2[e_g^c(\mathbf{F}) - i_g^c(\mathbf{F})]}$ . (Remember that  $g$  is in a safe forest if (II.3.44-46) apply).

Let  $g$  be a biped of a safe forest  $\mathbf{F}$ , and let us admit for a moment that inside  $g/A_{\mathbf{F}}(f)$  everything proceeds like in the preceding case. In the multiscale representation, this means that after internal spatial integrations are performed in  $g/A_{\mathbf{F}}(f)$  and internal line indices are summed with respect to the internal scale  $i_g^c(\mathbf{F})$  we obtain a quadratically divergent overall factor  $M^{2i_g^c(\mathbf{F})}$ , since  $\omega(g) = -2$  (apart from an overall  $K^{n(g)}$  and renormalon factors which will be considered later). The indices of the two external lines of  $g$  are  $e_g^c(\mathbf{F})$  and  $j$ , with  $j \geq e_g^c(\mathbf{F})$  by (II.3.35). We want to compare the two lines  $C^j(y-x)C^{e_g^c(\mathbf{F})}(x, z)$  with insertion of the biped counterterm for  $g$  at  $x$ , to a single ordinary line  $C^{e_g^c(\mathbf{F})}(y-z)$  (see Fig.II.3.7).

We have to add to the quadratically divergent factor  $M^{2i_g^c(\mathbf{F})}$  the effect of spatial integration over the position  $x$  of the  $\phi^2$  reduction vertex corresponding to  $g$ . This integration may be performed with the line of index  $j$  and brings a factor  $M^{-4j}$ . We combine this factor with the  $M^{2j}$  power counting factor of the line of index  $j$  and sum over  $j > e_g^c(\mathbf{F})$ ; this results in a factor  $M^{-2e_g^c(\mathbf{F})}$ . Using the triangular inequality we can also reconstruct the regular spatial decay of scale  $e_g^c(\mathbf{F})$  between  $y$  and  $z$ . Combining all these factors with formulae (II.3.44-46) we reconstruct the ordinary estimate (II.1.6) for  $C^{e_g^c(\mathbf{F})}(y-z)$ , plus a total power counting factor of  $M^{-2[e_g^c(\mathbf{F}) - i_g^c(\mathbf{F})]}$  for the mass counterterm and of 1 for the wave function counterterm. When we sum over  $i_g^c(\mathbf{F})$ , respecting the constraint of safeness for  $g$  which is  $i_g^c(\mathbf{F}) \leq e_g^c(\mathbf{F})$ , we get a constant in the first case and a factor  $e_g^c(\mathbf{F})$  in the second case (which may be bounded by  $i_{\max}(\mu)$ ).

The conclusion is that useless wave function counterterms (which correspond to marginal operators) generate the same kind of renormalon effects as useless coupling constant counterterms, and that useless mass counterterms do not. This is consistent with the bound (II.2.18) and the general definition of  $f(G)$ .

It remains to explain what we assumed, namely that inside each separate  $g/A_{\mathbf{F}}(g)$ ,  $g \in \mathbf{F} \cup \{G\}$  an analysis similar to the biped-free case takes place. The

key point is to check that after spatial integration, vertical index space decay is correctly generated. This means in particular that each connected component  $[g/A_{\mathbf{F}}(g)]_k^i$  with four or two external legs is properly renormalized, i.e. provides at each scale a reward at least  $M^{-1}$  instead of costing respectively a constant or  $M^2$ . Consider a proper closed subgraph  $h \in \mathbf{H}_\mu^c(\mathbf{F})$  such that  $B_{\mathbf{F}}(h) = g$ . The main remark is that the less stringent definitions (II.3.34-35) now allow renormalization of  $h$  on a wider range of indices than the minimal range in which  $h/A_{\mathbf{F}}(h)$  appears as a  $[g/A_{\mathbf{F}}(g)]_k^i$  component. Hence the bonuses  $M^{-1}$  or  $M^{-3}$  (respectively for  $h$  quadruped or biped) occur on this wider range, and this results in renormalized power counting not only for the corresponding proper closed graphs, but also for one-particle reducible bipeds and quadrupeds and for 1PI open quadrupeds when they occur as  $[g/A_{\mathbf{F}}(g)]_k^i$ 's. As expected, renormalization of proper closed subgraphs takes care of all others divergences. We will check this fact on typical examples.

Let us for simplicity forget about reduction by  $\mathbf{F}$ , i.e. write simply  $g$  for  $g/A_{\mathbf{F}}(g)$  etc... We consider first the case of a one particle reducible divergent subgraph  $h$  of  $g$ . There must exist a proper biped  $h' \subset h$  as in Fig.II.3.8, i.e.  $h'$  is a proper component of  $h$  and there is a common border vertex for  $h$  and  $h'$ ; otherwise  $h$  could not be divergent.

In this case consider the set  $I$  of indices  $i$  for which  $h$  is a connected component of  $g^i$ . We claim that if  $I$  is not empty,  $h' \in \mathbf{H}_\mu^c(\mathbf{F})$  and  $I$  is a piece of the extra range in which renormalization of  $h'$  provides an  $M^{-3}$  bonus. Indeed if  $i \in I$  one must have  $j_1 \geq i > j_2$  where  $j_1$  and  $j_2$  are the external indices of  $h'$  as shown in Fig.II.3.8. Therefore  $e_{h'}^c(\mathbf{F}) = j_2$ . Since  $i_{h'}(\mathbf{F}) \geq i$  (here we do mean  $i_{h'}(\mathbf{F})$  as defined in (II.3.4), *not*  $i_{h'}^c(\mathbf{F})$ ), we have  $i_{h'}^c(\mathbf{F}) \geq i_{h'}(\mathbf{F}) > j_2 = e_{h'}^c(\mathbf{F})$ , hence  $h' \in \mathbf{H}_\mu^c(\mathbf{F})$ . The factors earned in the renormalization of  $h'$  which cover the range  $]j_2, \min\{j_1, i_{h'}(\mathbf{F})\}]$ , in which obviously  $h'$  itself is not a connected component of  $g^i$ , can then be safely attributed to the renormalization of  $h$ ; they cover the desired set  $I$  of indices.

Finally we examine the case of an open quadruped  $h$ , for which the renormalization of its closure  $h^*$  should save the day (see Fig.II.3.5). Indeed by definition (II.3.34), if the set  $I$  of indices where  $h$  is a connected component of  $g^i$  is non empty we can conclude that  $h^*$  belongs to  $\mathbf{H}_\mu^c(\mathbf{F})$ ; and  $I$  is covered again by the index range  $i_{h^*}^c(\mathbf{F}) \geq i > i_{h^*}(\mathbf{F})$  in which we get bonus factors  $M^{-3}$  from the renormalization of  $h^*$ .

To complete the argument, let us check that there is no overlap, for a given biped  $h$ , between the factors used in the different cases considered in Fig.II.3.8 and II.3.5. Indeed the range  $i_h < i \leq i_h^c$  takes care of open quadrupeds inside  $h$ ; the range  $e_h < i \leq i_h$  takes care of  $h$  itself, and the range  $e_h^c < i \leq \min\{i_h, e_h\}$  takes care of one particle reducible graphs with  $h$  at one of their “ends”, and there is no overlap between these three ranges.

As a last remark, consider that when  $G$  itself is a biped, the derivatives corresponding to the  $(1 - \tau_G)$  operator apply to a true external leg of  $G$ ; in the formula analogous to (II.3.47) factors of degree 3 appear, which onez can bound by  $|p|^3$ . This explains the external momentum dependence of (II.2.18) and the definition of  $\hat{N}$  in this case. This completes our sketch of the general proof of Theorem II.2.1, and we refer to [dCR1] [FMRS2] for a somewhat more formalized

but less pedagogical proof.

As a technical remark, we note that the proof given above relies on the particular properties of  $\phi_4^4$  amplitudes and is therefore in some respects less general and systematic than the one of [FMRS2]. But in [FMRS2] a more inductive approach is used, in which a large number of Taylor operators may apply to the same line. This makes the inductive bound to be proven more complicated, and also results in a bound which is not uniform in  $n$  as far as the external momentum dependence is concerned, which for applications may be considered a significant drawback.

Again a byproduct of the analysis is that Theorem II.3.1 holds in the general case as well: usefully renormalized amplitudes  $I_G^{UR}$ , which are defined by the same formula (II.3.33) than in the biped free case, but with  $\mathbf{D}_\mu$  replaced by the generalized definition  $\mathbf{D}_\mu^c \equiv \mathbf{H}_\mu^c(\emptyset)$  still do not develop any factorial and satisfy the bound (II.3.33). Hence it is a general phenomenon that “useless counterterms” are solely responsible for renormalon effects.

In fact more precise statements can be derived from the analysis above. We remarked that mass counterterms do not contribute to any factorial effect; so there should be a way to write amplitudes with full mass renormalization and useful coupling constant and wave function renormalization which does not display any factorial behavior. It is easy to check (by trying some examples) that a proper biped can never overlap with any proper closed divergent subgraph, and this was proved rigorously in [dCR1]. In particular the set of all proper bipeds of  $G$  is itself a closed forest called  $\mathbf{B}(G)$ , and the product  $\prod_{b \in \mathbf{B}(G)} (1 - \tau_b)$ , factorizes in the  $\mathbf{R}$  operator. We are interested in factorizing only the mass renormalizations, so we write:

$$\mathbf{R} = \prod_{b \in \mathbf{B}(G)} (1 - \tau_b^{0*}) \sum_{\mathbf{F}} \prod_{b \in \mathbf{F}} (-\tau_b^{1*}) \prod_{q \in \mathbf{F}} (-\tau_q^*) \quad (\text{II.3.48})$$

where the first product runs over all proper bipeds  $b$  of  $G$ , the second over all bipeds of  $\mathbf{F}$  and the third over the quadrupeds of  $\mathbf{F}$ , and  $\tau_b^{0*}$  and  $\tau_b^{1*}$  are defined in (II.3.44-46). We may now apply the classification of forests only to the sum over forests in (II.3.48) and obtain amplitudes  $A_{G,\mathbf{F}}^{MR}$  in which the superscript  $MR$  means that all the mass renormalizations are fully performed. The case  $A_{G,\emptyset}^{MR} \equiv A_G^{MR,UR}$  is what we are looking for, and we obtain:

### Theorem II.3.2

$$|A_G^{MR,UR}(p_1, \dots, p_N)| = \left| \sum_{\mu} \int \prod_{v} dx_v \prod_{b \in \mathbf{B}(G)} (1 - \tau_b^{0*}) \prod_{b \in \mathbf{D}_\mu^c} (1 - \tau_b^{1*}) \prod_{q \in \mathbf{D}_\mu^c} (1 - \tau_q^*) \dots \right. \\ \left. \dots \prod_l C^{ii(\mu)}(x_l, y_l) \prod_{j=1}^N e^{ip_j x_j} \right| \leq (1 + \sup_j |p_j|^{\hat{N}}) K^{n(G)} \quad (\text{II.3.49})$$

Theorems II.3.1 and II.3.2 are powerful motivations to find the expansions corresponding to these usefully renormalized amplitudes which have the big advantage to be free of renormalon effects, hence to behave at large order like the amplitudes of a superrenormalizable theory. The next section is devoted to this problem.

From now on we drop the superscript  $c$  most of the time; in the biped free case the definitions of the first part of this section are to be used, and in the general case, the general definitions with superscript  $c$ .

## II.4 The effective expansion

*Tout ce qui est simple est inexact  
mais tout ce qui est compliqué est inutilisable.*  
— P. Valéry.

Let us summarize the conclusions of the preceding section. The bare expansion does not have a limit when the ultraviolet cutoff is removed. The renormalized expansion has cured this defect, but the price to pay for that is definitely too heavy. The “useless counterterms” make the study of renormalization quite painful (because one has to use forestry and to work always in reduced subgraphs  $g/A_{\mathbf{F}}(g)$ ); they also generate renormalon behavior which puts in danger the constructive program. Intuitively a factor  $K^n$  in large order estimates may be compensated by requiring the coupling constant to be very small, but this is not true for factorials (renormalon effects).

So we search for an expansion “in between” the bare and the renormalized expansion, as shown in Fig.II.4.1, one which is just a reshuffling of both, but with amplitudes which are the usefully renormalized ones. Hence it would have the advantage of ultraviolet finiteness without the renormalon effects. The idea of renormalization is that counterterms are hidden in the bare coupling constants; so we should get rid of the useless counterterms in this way. But since the useless counterterms depend on an index, or momentum scale, the corresponding coupling constants must also depend on it. Therefore the corresponding expansion cannot be a power series in a single bare or renormalized constant, but must be a series in a whole sequence of index dependent constants called the effective constants. This is particularly clear on equations (II.4.1-2) below: the separation of counterterms into useful and useless pieces being index dependent, there should be a compensating index dependence in the effective constants in order for the left hand side of (II.4.1-2) to be scale independent.

We pause briefly to remark that we have reached the typical renormalization group concept of effective or running constants from the unusual point of view of simply organizing counterterms so as to get the best large order bounds. The historical and more traditional road to these concepts has been to write renormalization group equations which investigate the dependence of the theory upon the subtraction scale, or the somewhat equivalent problem of finite changes of the counterterms [SP][GL]. This standard road paved the way for an enormous amount of work, from the discovery of the particularly interesting Callan-Symanzik equation [Ca1][Sy1], to Wilson’s generalization of the renormalization group program [Wil][KW], the investigation of renormalization group fixed points and their stability for various models, and led even by a back-reaction to new proofs of renormalizability [Ca2]. We do not try any review of this vast subject. In particular we do not enter into the definition of normal ordered products, the Zimmermann identities, the renormalization group equations which investigate how the theory changes with a change of cutoff or a change of the subtraction scale. This is because these topics and the Callan-Symanzik equation are explained in great detail in most books devoted to quantum field theory or renormalization theory: a ded-

icated reference is [Co]. For a rather complete overview of the renormalization group and its application to critical phenomena, we refer to [Am]. The reader may be disappointed not to find this standard material here; but the traditional presentation involves defining the corresponding concepts and equations as bare or renormalized power series, and we are about to argue that the proper framework is neither the bare nor the renormalized expansions but an effective expansion in between. We are also going to define the discrete version of these equations which is naturally adapted to the multiscale representation. In particular we will find natural discrete analogue of the famous renormalization group functions such as the  $\beta$  function, and we prefer to proceed directly toward this goal. Discrete equations are less elegant than differential equations, but again they seem to be required by the constructive point of view to be developed in part III.

However let us warn the reader that the effective expansion derived below is more limited in scope than the general philosophy of the renormalization group, and also stress again that the multiscale slicing of the propagator is a somewhat simpler but less general method than Kadanoff's and Wilson's block-spin methods; as remaked already, it requires indeed a gaussian measure in the problem. We hope however that the presentation below may have an advantage over the more standard one, at least for the beginner: it applies to the internal lines "inside" the Feynman amplitudes rather than to the external legs. Hence it develops the correct intuition that renormalization group behavior is an essential piece for the *construction* of the theory, not just a device to analyze its behavior in various regimes.

Let us make precise the paths by which the three expansions of Fig.II.4.1 communicate. The bare expansion leads to the effective one by developing selectively some of the bare constants into effective constants plus the useful piece of the counterterms. This is the path used in chapter III, where the bare theory is always the constructive starting point. Conversely the renormalized expansion leads to the effective one by resumming the useless counterterms, or equivalently by absorbing them into effective constants. Hence in a schematic way:

$$\text{bare constant} = \text{effective constants} + \text{useful counterterms} \quad (\text{II.4.1})$$

$$\text{renormalized constant} = \text{effective constants} - \text{useless counterterms} \quad (\text{II.4.2})$$

these equations being consistent with the standard rule:

$$\text{bare constant} = \text{renormalized constant} + \text{full counterterms} \quad (\text{II.4.3})$$

The two paths are not exactly symmetrical, however, and one can consistently argue that the first one is shorter and simpler than the second one. This is because (forgetting for the moment the subtleties due to bipeds) the definitions (II.1.10-11)-(II.3.4-5) of indices  $i_g$  and  $e_g$  make the condition of almost locality  $i_g > e_g$  rather stringent: *every* internal line of  $g$  has to have higher index than *every* external line of  $g$ . In other words there is not much renormalization performed in the usefully renormalized expansion and in this way it is closer to the bare one than to the renormalized one. Perhaps this argument may not seem serious at first sight. Nevertheless it is also the reason for which the first path is simpler; we

do not expect any sum over forests to appear along it, since these sums neither appear in the formulas for the starting point (the bare amplitudes) nor for the end point (the usefully renormalized amplitudes, see (II.3.33), and also (II.3.49)). Sums over various forests in fact only arise when counterterms are developed for intermediate situations in which internal and external indices are quite mixed.

It is therefore natural to start with a theorem relating the bare and effective expansion. As in the preceding section and for the same reasons we state it first for the biped-free piece of the perturbative expansion, then extend it to the general case.

We fix a cutoff index  $\rho$  and the bare coupling  $g_\rho$ . For each vertex  $v$  of a graph  $G$  it is useful to define;

$$e_v(\mu) = \max\{i_l(\mu) | l \text{ hooked to } v\} \quad (\text{II.4.4})$$

Recall that by convention the index of external lines is  $-1$ , so  $e_v(\mu) = -1$  is possible but only for the unique vertex of the trivial graph with a single vertex,  $N = 4$ , and no internal lines.

The bare expansion for a connected Schwinger function with cutoff  $\rho$  is written in analogy with (I.4.15) as:

$$C_{N,bf}^\rho = \sum_{G,\mu | \mu \leq \rho} \frac{(-g_\rho)^{n(G)}}{S(G)} A_{G,\mu} \quad (\text{II.4.5})$$

where the sum is over assignments  $\mu \in [0, \rho]^{l(G)}$  (in short  $\mu \leq \rho$ ) and over connected biped free graphs with  $N(G) = N$ , as indicated by the index  $bf$  (biped-free). (II.4.5) defines  $C_{N,bf}^\rho$  as a formal power series in  $g_\rho$ .

#### **Theorem II.4.1: Existence of the effective expansion**

There exist  $\rho + 1$  formal power series in  $g_\rho \equiv g_\rho^\rho$ , called  $g_{\rho-1}^\rho, g_{\rho-2}^\rho, \dots, g_0^\rho$  and  $g_{-1}^\rho$  (the upper index is to remind the reader that the entire theory has ultraviolet cutoff  $\rho$ ) such that the formal power series (II.4.5) is the same as:

$$C_{N,bf}^\rho = \sum_{G,\mu | \mu \leq \rho} \left[ \prod_{v \in G} (-g_{e_v(\mu)}^\rho) \right] \frac{1}{S(G)} A_{G,\mu}^{UR} \quad (\text{II.4.6})$$

where we recall (see (II.3.33)):

$$A_{G,\mu}^{UR} \equiv \int \prod_v dx_v \prod_{h \in \mathbf{D}_\mu} (1 - \tau_h^*) Z_{G,\mu} \quad (\text{II.4.7})$$

$Z_{G,\mu}$  is the integrand for graph  $G$  and assignment  $\mu$ ; for instance if we take the external arguments to be fixed momenta we have:

$$Z_{G,\mu} \equiv \prod_l C^{i_l(\mu)}(x_l, y_l) \prod_{j=1}^N e^{ip_j x_j} \quad (\text{II.4.8})$$

and the effective constants  $g_i^\rho$  obey the following inductive definition:

$$g_i^\rho = g_{i+1}^\rho - \sum_{\substack{H \text{ quadruped}, \mu \leq \rho \\ i_H(\emptyset) = i+1}} \frac{1}{S(H)} \prod_{v \in H} (-g_{e_v(\mu)}^\rho) \int \prod_{v \in H} dx_v$$

$$\left[ \prod_{h \in \mathbf{D}_\mu(H), h \neq H} (1 - \tau_h^*) \right] \tau_H Z_{H,\mu} \quad (\text{II.4.9})$$

where  $\tau_H Z_{H,\mu} = Z_{H,\mu}|_{p=0} = \prod_{l \in H} C^{i_l(\mu)}(x_l, y_l)$  is the integrand for  $H$  taken at 0 external momenta, and  $\mathbf{D}_\mu(H)$  is simply the forest  $\mathbf{D}_\mu = \mathbf{H}_\mu(\emptyset)$  of the preceding section, but for the graph  $H$ . In (II.4.9), the summation over quadrupeds does not include the trivial case of the graph reduced to a single vertex, which corresponds in fact to the first factor  $g_{i+1}^\rho$  in the right hand side of (II.4.9).

The minus sign in (II.4.9) is consistent with the minus sign in (II.4.2), because the counterterm is  $\int [ \prod_{h \in \mathbf{D}_\mu(H), h \neq H} (1 - \tau_h^*) ] (-\tau_H) Z_H$ , but the vertex has really a value  $-g$ , so the reader should think of (II.4.9) as a more convenient form of the equation  $-g_i = -g_{i+1} - \sum_H \text{counterterm}(H)$ .

(II.4.9) defines each  $g_i^\rho$  (by inductive substitution) as a formal power series in  $g_\rho$  of the form  $g_\rho + \sum_{n \geq 2} \gamma_n^i(g_\rho)^n$ . The induction stops at  $g_{-1}^\rho$  which is the last one for which the sum in (II.4.9) is not empty. Let us apply the result (II.4.6) to  $N = 4$  and put to 0 the four external momenta. When  $G$  is a non trivial quadruped,  $G$  itself always belongs to  $\mathbf{D}_\mu(G)$ , and the  $(1 - \tau_G)$  operator makes  $A_{G,\mu}^{UR}$  vanish at 0 external momenta. For the trivial graph with a single vertex  $v$  we remarked that  $e_v(\mu) = -1$ . Hence the formal power series in  $g_\rho$  (II.4.6) for  $C_4^\rho(0,0,0,0)$  reduces exactly to  $-g_{-1}^\rho$ . This means that in the sense of formal power series in  $g_\rho$  we must identify  $g_{-1}^\rho$  with the renormalized coupling  $g_r$ , which by definition of our subtraction scheme is precisely minus the connected four point function at 0 external momenta  $C_4^\rho(0,0,0,0)$ .\*

The proof of Theorem II.4.1 is a simple combinatoric exercise; no analysis is involved, since all integrals involved have cutoffs and are therefore obviously absolutely convergent. Again the combinatoric has to be checked at the level of contraction schemes. We go from (II.4.5) to (II.4.6) by pulling out inductively the useful counterterms hidden in  $g_\rho$ , one slice after the other. At slice  $i$  an intermediate version of Theorem II.4.1 is obtained:

$$C_{N,bf}^\rho = \sum_{G,\mu \leq \rho} \left[ \prod_{v \in G} (-g_{\sup(i, e_v(\mu))}^\rho) \right] \frac{1}{S(G)} A_{G,\mu}^{UR,i} \quad (\text{II.4.10})$$

where:

$$A_{G,\mu}^{UR,i} \equiv \int \prod_v dx_v \prod_{h \in \mathbf{D}_\mu^i} (1 - \tau_h^*) Z_{G,\mu} \quad (\text{II.4.11})$$

and:

$$\mathbf{D}_\mu^i \equiv \{h \in \mathbf{D}_\mu | i_h > i\} \quad (\text{II.4.12})$$

(II.4.10) is obviously nothing but (II.4.5) if  $i = \rho$ . Assuming it at scale  $i + 1$ , we prove it at scale  $i$  by simply adding and subtracting the counterterms which change  $A_{G,\mu}^{UR,i+1}$  into  $A_{G,\mu}^{UR,i}$ . These are the counterterms corresponding to the

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\* The renormalization condition of BPHZ define in fact  $g_r$  as the 0 momentum value of  $-[\frac{\partial}{\partial p^2} \Gamma_2(0)]^{-2} \Gamma_4(0)$ , hence include a field strength renormalization factor which disappears only if the renormalized parameter  $a_r$  is 1. However we do not need to discuss this subtlety for the moment, since in the biped-free theory there is obviously no wave function renormalization.

quadrupeds  $\{H_1, \dots, H_k\} = \{H \in \mathbf{D}_\mu | i_H = i + 1\}$ . Hence we add and subtract to each  $A_{G,\mu}^{UR,i+1}$  the quantity:

$$\sum_{\substack{S \subseteq \{H_1, \dots, H_k\} \\ S \neq \emptyset}} \prod_{H_j \in S} (-\tau_{H_j}) \prod_{h \in \mathbf{D}_\mu^{i+1}} (1 - \tau_h^*) Z_{G,\mu} \quad (\text{II.4.13})$$

The piece added changes  $\prod_{h \in \mathbf{D}_\mu^{i+1}} (1 - \tau_h^*)$  into  $\prod_{h \in \mathbf{D}_\mu^i} (1 - \tau_h^*)$  in each amplitude, hence it changes  $A_{G,\mu}^{UR,i+1}$  into  $A_{G,\mu}^{UR,i}$ . The piece “subtracted” should be developed as a sum over  $S$ , so as to get:

$$C_{N,bf}^\rho = \sum_{\substack{(G,\mu,S), \mu \leq \rho \\ S \subseteq \mathbf{D}_\mu^i - \mathbf{D}_\mu^{i+1}}} [\prod_{v \in G} (-g_{\sup(i+1, e_v(\mu))}^\rho)] \frac{1}{S(G)} A_{G,\mu,S}^{UR,i} \quad (\text{II.4.14})$$

with

$$A_{G,\mu,S}^{UR,i} \equiv A_{G,\mu}^{UR,i} \text{ if } S = \emptyset \quad (\text{II.4.15})$$

and

$$A_{G,\mu,S}^{UR,i} \equiv \int \prod_v dx_v \prod_{H_j \in S} (-\tau_{H_j}) \prod_{h \in \mathbf{D}_\mu^{i+1}} (1 - \tau_h^*) Z_{G,\mu} \text{ otherwise.} \quad (\text{II.4.16})$$

Remark that this induction is really nothing but the Bogoliubov induction, but with the additional element that the multiscale decomposition provides at each scale a well defined natural family  $H_1, \dots, H_k$  of disjoint subgraphs, which are the ones to which the Bogoliubov induction should be applied, instead of being performed blindly with respect to scales.

We can now define, since the elements of  $S$  are disjoint, the collapse  $\phi_i$  as an operation which is defined on triplets  $(G, \mu, S)$ ,  $S \subseteq \mathbf{D}_\mu^i - \mathbf{D}_\mu^{i+1}$ , and which sends  $(G, \mu, S)$  to  $(G', \mu', \emptyset)$ ,  $G'$  being obtained from  $G$  by reducing each  $H_j \in S$  to a single vertex, and  $\mu'$  being the assignment derived from  $\mu$  by simple restriction to the lines of  $G'$ . Remark that every vertex of  $G'$  corresponding to such a reduction must have  $e_v(\mu) = e_v(\mu') \leq i$ . We reorder now (II.4.14) as:

$$C_{N,bf}^\rho = \sum_{(G', \mu')} \{ \sum_{\substack{(G,\mu,S), \mu \leq \rho \\ \phi_i(G,\mu,S) = (G', \mu', \emptyset)}} [\prod_{v \in G} (-g_{\sup(i+1, e_v(\mu))}^\rho)] \frac{1}{S(G)} A_{G,\mu,S}^{UR,i} \} \quad (\text{II.4.17})$$

For each  $(G', \mu')$  the corresponding sum in (II.4.17) is an infinite power series which in fact replaces exactly, at each vertex  $v$  of  $G'$  satisfying  $e_v(\mu') \leq i$ , the coupling  $g_{i+1}^\rho$  by the right hand side of (II.4.9), hence by  $g_i^\rho$ ; the sum over  $H$  in (II.4.9) indeed corresponds exactly to the sum over all possible insertions of an  $H_j$  which is collapsed by the  $\phi_i$  operation to the vertex  $v$ , in the above notation. To check that combinatoric factors agree, one has again to perform this analysis at the level of contraction schemes, in which case it becomes the same problem as for the Bogoliubov induction considered above.

This achieves the proof of (II.4.10) at scale  $i$ , hence by induction, the proof of Theorem II.4.1.

It remains to show that the effective expansion has the advantages but not the drawbacks of the renormalized expansion, and so we should derive a version in which the ultraviolet limit has been taken. This is not straightforward, since we can no longer use the bare constant. But since  $g_r = g_{-1}^\rho$  is a formal power series in  $g_\rho$  starting with  $g_\rho$ , it is possible to invert it at the level of formal power series. This is also true for each effective constant  $g_i^\rho$  which by substitution becomes a formal power series in  $g_r$  starting with  $g_r$ . (Remark that the power series  $g_i^\rho$  obtained in this way still depend on  $\rho$ , the global cutoff).

Now the series (II.4.6) when considered as formal power series in  $g_r$  through such substitutions is nothing but exactly the ordinary (biped-free) fully renormalized power series in  $g_r$  with cutoff  $\rho$  (i.e. with full renormalization operator  $\mathbf{R}$ , but propagators  $C^\rho$  instead of  $C$ ).

An ultraviolet limit of the effective expansion may then be obtained in the following sense:

#### **Theorem II.4.2 Ultraviolet limit of the effective expansion**

The effective constants  $g_i^\rho$  have a limit as  $\rho \rightarrow \infty$  order by order as formal series in  $g_r$ . This limit is called  $g_i^\infty$ , or simply  $g_i$ . Furthermore, in the sense of formal power series in  $g_r$ , the biped free part of the BPHZ renormalized expansion for  $C_N$  satisfies:

$$C_{N,bf} = \sum_{G,\mu} \left[ \prod_{v \in G} (-g_{e_v(\mu)}) \right] \frac{1}{S(G)} A_{G,\mu}^{UR} \quad (\text{II.4.18})$$

The subtle point is to show the convergence (order by order in  $g_r$ ) of  $g_i^\rho$  to  $g_i$ . This may be done by rewriting  $g_i^{\rho+1} - g_i^\rho$  as a sum over renormalized graphs which have at least one line at scale  $\rho$ , and then use the vertical exponential decay of such graphs in  $\rho - i$ . Formula (II.4.18) then follows from the similar statement for the theory with cutoff and the absolute convergence of the renormalized amplitudes with cutoff to the ones without cutoff.

Theorem II.4.2 achieves our goal of an effective expansion which is ultraviolet finite, like the renormalized one, but free of renormalons and of sums over forests.

When bipeds are added, we may choose between several generalizations of Theorems II.4.1-II.4.2. The most natural generalization is to derive an effective expansion with three types of effective parameters, the effective coupling constant, the effective mass and the effective wave function constant. This requires to re-express that theory as a sum over generalized  $\phi^4$  graphs with two point mass and wave function insertions, as in Fig.I.4.2 (recall that the initial mass and wave function parameters are not expressed as coupling constants, but used to build the propagator of the theory).

To implement this idea it is convenient to use the general definitions of Sect.II.3B for internal and external indices. We have now generalized graphs  $\hat{G}$  with regular vertices  $v \in V(\hat{G})$  with 4 legs, for which definition (II.4.4) is adequate, and two other sets of vertices  $W^0(\hat{G})$  and  $W^1(\hat{G})$  with 2 legs, corresponding respectively to mass and wave function insertions. For  $w \in W^0 \cup W^1$ , in view of (II.3.35) one should define

$$e_w(\mu) = \min\{i_l(\mu) | l \text{ hooked to } w\} \quad (\text{II.4.19})$$

Then we have the generalization of Theorem II.4.1:

**Theorem II.4.3**

There exist  $3(\rho + 1)$  formal power series in  $g_\rho$ , called  $g_i$ ,  $\delta m_i^2$  and  $\delta a_i$ ,  $i = \rho - 1, \dots, 0, -1$  (they depend on  $\rho$ , like those of Theorem II.4.1, but we drop this dependence to avoid too heavy notations), such that the formal power series in  $g_\rho$  for  $C_N^\rho$  can be rewritten as:

$$\begin{aligned} C_N^\rho = & \sum_{\hat{G}, \mu} \left[ \prod_{v \in V(\hat{G})} (-g_{e_v(\mu)}) \right] \left[ \prod_{w \in W^0(\hat{G})} (-\delta m_{e_w(\mu)}^2) \right] \\ & \left[ \prod_{w \in W^1(\hat{G})} (-\delta a_{e_w(\mu)}) \right] \frac{1}{S(G)} A_{\hat{G}, \mu}^{UR} \end{aligned} \quad (\text{II.4.20})$$

where the formula for  $A_{\hat{G}, \mu}^{UR}$  is:

$$A_{\hat{G}, \mu}^{UR} \equiv \int \prod_{v \in V \cup W^0 \cup W^1} dx_v \prod_{h \in \mathbf{D}_\mu} (1 - \tau_h^*) \prod_{w \in W^1(\hat{G})} (-\Delta) Z_{\hat{G}, \mu} \quad (\text{II.4.21})$$

$Z_{\hat{G}, \mu}$  is the integrand for  $\hat{G}$ , and  $\mu$ , namely at fixed external momenta:

$$Z_{\hat{G}, \mu} \equiv \prod_l C^{i_l(\mu)}(x_l, y_l) \prod_{j=1}^N e^{ip_j x_j} \quad (\text{II.4.22})$$

and the operator  $\Delta \equiv \partial_\nu \partial_\nu$  acts, for each  $w \in W^1(\hat{G})$ , on one of the two propagators hooked to  $w$ , in agreement with formula (II.3.46) for the action of the operator  $\tau^{1*}$ .

Furthermore the effective constants  $g_i$ ,  $\delta m_i^2$ ,  $\delta a_i$  obey inductive relations generalizing (II.4.9). One starts with  $g_\rho$  being the bare coupling,  $\delta m_\rho^2 \equiv 0$  and  $\delta a_\rho \equiv 0$  and define by recursion:

$$\begin{aligned} g_i = g_{i+1} - & \sum_{\substack{\hat{H} \text{ quadruped}, \mu \leq \rho \\ i_{\hat{H}}(\emptyset) = i+1}} \frac{1}{S(\hat{H})} \prod_{v \in V(\hat{H})} (-g_{e_v(\mu)}) \left[ \prod_{w \in W^0(\hat{H})} (-\delta m_{e_w(\mu)}^2) \right] \\ & \left[ \prod_{w \in W^1(\hat{H})} (-\delta a_{e_w(\mu)}) \right] \int \prod_{v \in V \cup W^0 \cup W^1(\hat{H})} dx_v \left[ \prod_{h \in \mathbf{D}(\hat{H})_\mu, h \neq \hat{H}} (1 - \tau_h^*) \right] \tau_{\hat{H}} Z_{\hat{H}, \mu} \end{aligned} \quad (\text{II.4.23})$$

where  $\tau_{\hat{H}} Z_{\hat{H}, \mu} = Z_{\hat{H}, \mu}|_{p=0} = \prod_{l \in H} C^{i_l(\mu)}(x_l, y_l)$  is the integrand for  $\hat{H}$  taken at 0 external momenta;

$$\begin{aligned} \delta m_i^2 = \delta m_{i+1}^2 + & \sum_{\substack{\hat{B} \text{ biped}, \mu \leq \rho \\ i_{\hat{B}}(\emptyset) = i+1}} \frac{1}{S(\hat{B})} \prod_{v \in V(\hat{B})} (-g_{e_v(\mu)}) \left[ \prod_{w \in W^0(\hat{B})} (-\delta m_{e_w(\mu)}^2) \right] \\ & \left[ \prod_{w \in W^1(\hat{B})} (-\delta a_{e_w(\mu)}) \right] \int \prod_{v \in V \cup W^0 \cup W^1(\hat{B})} dx_v \left[ \prod_{h \in \mathbf{D}(\hat{B})_\mu, h \neq \hat{B}} (1 - \tau_h^*) \right] Z_{\hat{B}, \mu}^0 \end{aligned} \quad (\text{II.4.24})$$

$$\delta a_i = \delta a_{i+1} + \sum_{\substack{\hat{B} \text{ biped}, \mu \leq \rho \\ i_{\hat{B}}(\emptyset) = i+1}} \frac{1}{S(\hat{B})} \prod_{v \in V(\hat{B})} (-g_{e_v(\mu)}) \left[ \prod_{w \in W^0(\hat{B})} (-\delta m_{e_w(\mu)}^2) \right] \\ \left[ \prod_{w \in W^1(\hat{B})} (-\delta a_{e_w(\mu)}) \right] \int \prod_{v \in V \cup W^0 \cup W^1(\hat{B})} dx_v \left[ \prod_{h \in D(\hat{B})_\mu, h \neq \hat{B}} (1 - \tau_h^*) \right] Z_{\hat{B}, \mu}^1 \quad (\text{II.4.25})$$

where  $Z_{\hat{B}, \mu}^0 = \prod_{l \in H} C^{i_l(\mu)}(x_l, y_l)$  and  $Z_{\hat{B}, \mu}^1 = \frac{1}{8}|x - y|^2 \prod_{l \in H} C^{i_l(\mu)}(x_l, y_l)$ ,  $x$  and  $y$  being the border vertices of  $\hat{B}$ . These last definitions are again consistent with the action of the Taylor operator for a biped as shown in (II.3.44-46).

This theorem is checked by induction exactly as the previous one. A similar theorem for the 1PI or vertex functions  $\Gamma_N$  is left to the reader to formulate, the sums being of course restricted everywhere in this case to one particle irreducible graphs.

Inverting the series  $g_i$  we can reexpress them as series (still depending on  $\rho$ ) in  $g_{-1} = g_r$  and substituting them everywhere we obtain the BPHZ renormalized expansion for  $C_N^\rho$ . In particular we can obtain renormalized masses and wave function constants  $m_r^2 = m^2 + \delta m_{-1}^2$  and  $a_r = 1 + \delta a_{-1}$ , if  $m^2$  and 1 were the bare values used to build the propagator. These renormalized parameters coincide with the ones of the BPHZ prescription at 0 external momenta, so that  $m_r^2 = -\Gamma_2(0)$  and  $a_r = -\frac{d}{dp^2}\Gamma_2(0)$ .

For these series we can pass to the ultraviolet limit  $\rho \rightarrow \infty$  and the generalization of Theorem II.4.2 also holds:

#### Theorem II.4.4

The effective expansion reexpressed in terms of  $g_r$ ,  $m_r$  and  $a_r$  in the limit  $\rho \rightarrow \infty$  is the same order by order in  $g_r$  as the usual BPHZ renormalized series.

It is in principle possible to work out a direct proof of Theorem II.4.2 and II.4.4, i.e. to reshuffle directly the renormalized series into effective ones without ever using an ultraviolet cutoff and the bare expansion. The effective constants come from the resummation of the useless pieces of the counterterms of the renormalized expansion. This approach has been partly implemented in [Ri1] (see also [DFR]), in which such explicit resummations have been introduced for all counterterms of the bubble type (the so called parquet forests). A general resummation rule in the sense of formal power series obviously exists by Theorem II.4.4, but to write it in the form of a simple set of explicit resummation rules seems difficult since forests are now an essential ingredient. This illustrates the fact that path 2 in Fig.II.4.1 is less straightforward than path 1.

Expansion (II.4.20) is not as canonical as (II.4.6). In particular it is rather natural to work out similar versions in which chains of two point insertions in the generalized graphs of (II.4.20) have been resummed. This is possible essentially because two point insertions lead to geometric, explicitly summable series. This leads to an effective expansion with effective coupling constants and effective *propagators*; the advantage is that only ordinary graphs are required, the drawback is that the formula for the effective propagators is somewhat complicated. We refer to [FMRS4] for an example of such a construction.

Finally one may take advantage of the fact that mass renormalization does neither create forest nor renormalons problems, as was shown in Theorem II.3.2;

hence one may derive a perturbative expansion which has full mass renormalization and effective coupling and wave function constants; the amplitudes for this expansion are the  $A_G^{MR,UR}$  of (II.3.49).

Furthermore in asymptotically free theories which will be discussed soon, it is even useful to consider an expansion with effective coupling constants, bare wave function constant and renormalized mass, since we will see that asymptotic freedom makes in fact the wave function renormalization finite. The corresponding amplitudes are noted  $A_G^{MR,CCUR}$  for “mass renormalized, coupling constant usefully renormalized” (and wave function not renormalized...). Such expansions, although of mixed character, are optimal for instance for the construction of the planar  $-g\phi_4^4$  theory which is the goal of the next chapter. They can also be used in the constructive context [FMRS4][FMRS5].

To summarize this section, the multiscale decomposition provides a simple, well defined prescription for writing down discrete flows for the relevant and marginal couplings of a renormalizable theory like  $\phi_4^4$ ; the bare or renormalized parameters simply provide particular boundary conditions for these flows. This formalism is completely rigorous at the level of formal power series, and it leads to effective expansions which are formally equivalent to the bare or renormalized ones but are both ultraviolet finite and free of forest complications and of renormalon effects.

## II.5 Construction of “wrong sign” planar $\phi_4^4$ .

In this section we show how to apply the formalism of the effective expansion beyond the sterile level of formal power series, on a simple model which is not a full-fledged field theory but has nevertheless some physical interest.

Even in the effective version without renormalon effects, it is not easy to sum up perturbation theory because of the large number of graphs involved at large order. This divergence of perturbation theory occurs even in 0 dimension, for a single integral  $\int e^{-x^2-gx^4} dx$ . In the next section we will analyze it in some detail and conclude that we must trade absolute summation for Borel summation at best. However there is a second problem. It is usual in the construction of a theory in the weak coupling regime to take care of combinatoric factors like the constants in the uniform theorems of Sect.II.1-2-3 by requiring the coupling constant to be small enough. But in the effective expansion of Sect.II.4 there is an infinite set of such coupling constants, and we need therefore them to be all simultaneously small enough. If perturbation theory is asymptotic, the sign of the first term in the recursion relation (II.4.9) or (II.4.23) for  $g_i - g_{i+1}$  becomes very important. If this sign is negative,  $g_{i+1}$  will be larger than  $g_i$  (for small enough  $g_i$ ) and holding the renormalized constant  $g_r = g_{-1}$  fixed to a small value it is very doubtful that all other effective constants  $g_i$  can be made simultaneously small. On the other hand if it is positive,  $g_{i+1}$  is smaller than  $g_i$  and the goal seems attainable. In this second case the theory is called asymptotically free (or more precisely ultraviolet asymptotically free); indeed in the cases where we can in fact make sense out of the recursion relation (II.4.9) beyond formal power series,  $g_i$  tends to 0 as  $i \rightarrow \infty$ , with  $g_r$  kept fixed and small.

The first term in (II.4.9) or (II.4.23) corresponds to  $H$  (or  $\hat{H}$ ) being the bubble graph. In the  $e^{-g} \int \phi^4$  theory with the ordinary sign of the coupling constant ( $g > 0$ ),  $\tau_H Z_H$  is positive; because of the minus sign in the exponential there is however a minus sign in (II.4.9)-(II.4.23) and the theory is not ultraviolet asymptotically free. This fact of life is mathematically (and perhaps even physically) frustrating because it deprives us of the simplest and most natural model to construct. However by reversing the direction in the index space, a theory has to be either asymptotically free in the ultraviolet or in the infrared direction; the ordinary  $\phi^4$  theory is therefore asymptotically free in the infrared regime, which opens up the possibility, exploited in part III, to construct the critical or massless limit of the model with fixed ultraviolet cutoff. An other idea is to change the sign of the coupling constant: we call the corresponding theory, with action  $e^{-g} \int \phi^4$  and  $g < 0$  the “wrong sign”  $\phi^4$  (or the “negative coupling”  $\phi^4$  theory, but this may lead to some confusion, because  $-g$  is positive for  $g < 0$ ...). This model is ultraviolet asymptotically free, but a negative  $\phi^4$  potential means that  $\phi = 0$  is an unstable minimum, and in spite of many efforts there is therefore no construction of the model (apart from analytic continuations which are expected not to meet the axiomatic requirements [GK5]).

But there is a restricted version of the wrong sign  $\phi^4$ , the planar wrong sign  $\phi^4$ , which offers a nice benchmark to study asymptotic freedom and make sense of the flows and recursion relations of the previous section beyond formal power series. The planar theory was proposed initially by 't Hooft [tH2] as an approximation for

the study of SU(N) gauge theories at large N; it has also been studied in connection with random surface models or string theory. The name planar comes from the fact that the perturbation theory of these models is roughly speaking restricted to graphs which can be drawn on a plane without self-intersections. Constructions of the wrong sign planar  $\phi^4$  were performed in [tH5-6-7], [Ri1], [GaNi]; here we try a compromise retaining some of the best aspects of the initial constructions. Because it is not adapted to the constructive versions of part III we do not however retain the nice continuous slicing provided by  $\alpha$  parameters and used in [Ri1][Hu]; we invite the reader in search of a “truly optimal” version to write down the necessary modifications as an exercise.

The  $\phi^4$  planar model which we consider is the  $N \rightarrow \infty$  limit of an  $N$  by  $N$  matrix valued  $\phi^4$  model. Hence the field  $\phi$  is a bosonic field with components  $\phi_{p,q}$ ,  $1 \leq p, q \leq N$ . The formal functional measure for this model is similar to (I.3.1) (for a real valued  $\phi$ ):

$$d\nu = \frac{1}{Z} e^{+(g/4!N) \int \text{Tr}(\phi^t \phi)^2 - (m^2/2) \int \text{Tr} \phi^t \phi - (a/2) \int \text{Tr} \partial_\mu \phi^t \partial^\mu \phi} D\phi \quad (\text{II.5.1})$$

where  $D\phi$  is a product of independent formal Lebesgue measures for each component of  $\phi$ . This model has global  $O(N)$  invariance (there is a complex version with global U(N) invariance). The “wrong sign” case in (II.5.1) corresponds to  $g > 0$ , which avoids minus signs in many of the formulas below.

The Feynman rules of this model are discussed for instance in [GrK]. Propagators carry a matrix index, hence it is convenient to represent them as double lines, one for each matrix index. They are shown together with vertices and some graphs in Fig.II.5.1. Remark the cyclic symmetry of the vertex. To each closed loop there corresponds a sum over possible values of the index flowing through this loop, which gives a corresponding factor  $N$ . There is a factor  $N^{-1}$  per vertex, hence the overall factor for a graph with  $k$  external lines is

$$N^{2-k-2h} \quad (\text{II.5.2})$$

where  $h$  is the number of handles of the surface on which the graph is a triangulation. This number may be visualized by considering each propagator (double line) as a thin ribbon. All the ribbons of the external legs are tied to a single point at infinity (without tangling them). Then each closed index loop is filled with a flat piece of surface matching with each ribbon on its boundary. In this way a compact surface is generated, and  $h$  is simply its genus (its number of holes or of “handles”). The formula (II.5.2) is then a standard homological formula, which can be checked on the examples of Fig.II.5.1.

At fixed  $k$  in the limit  $N \rightarrow \infty$  the leading term surviving in (II.5.2) corresponds to the sum of all graphs with no handles ( $h = 0$ ), and is called the planar theory. Apart from the constant overall factor  $N^{2-k}$  in front, the amplitudes are then exactly similar to the ones of the ordinary  $\phi^4$  theory, the only difference lying in the planar restriction. Remark that it might be misleading to simply state that the planar theory is the ordinary  $\phi^4$  theory reduced to “the graphs that can be drawn on a plane without self intersection”, because in the prescription for computing the value of  $h$  above we cannot avoid to consider the double lines. It is true that for any planar graph the collapse of the double propagators to single lines

gives an ordinary  $\phi^4$  graph that can be drawn on a plane without self crossings. But the combinatoric factors associated to such a graph are usually different in the planar and the ordinary theory; this is already true for the “bubble graph” (and leads for instance to different values of the first term in the  $\beta$  function for the ordinary and the planar theory).

From (II.5.2) only the two point function ( $k = 2$ ) has a non zero limit, so that the planar theory strictly speaking is a free field theory with a complicated propagator. In fact we will give a non trivial meaning to the sum over planar graphs for *any* Schwinger function  $S_k$ , discarding the constant overall factor  $N^{2-k}$  in front. This point of view allows to introduce the renormalized coupling constant as ususally as the value of the connected four point function at 0 external momenta, etc... Nevertheless one should keep in mind that these planar series without the overall factor  $N^{2-k}$  are no longer exactly the  $N \rightarrow \infty$  limit of (II.5.1), and in particular one should not believe that they correspond to a full fledged interacting field theory satisfying the Osterwalder-Schrader axioms.

There is an even simpler model with similar features: it is the  $N \rightarrow \infty$  limit of the  $N$ -vector  $\phi^4$  theory, which we pause to discuss briefly. In this model the action is similar to (II.5.1) but  $\phi$  is an  $N$  component vector, so that  $\phi^t\phi$  is a scalar product, and there is no need for traces in (II.5.1). The Feynman rules and typical graphs for this model are shown in Fig.II.5.2. The leading terms as  $N \rightarrow \infty$  for vacuum graphs, two, and four point functions behave as  $N$ , 1, and  $N^{-1}$ , respectively (some of them are pictured in Fig. II.5.2). Again this limit should be considered a free field. By Wick ordering one may eliminate the tadpoles (lines with both ends at the same vertex). Then the leading graphs for instance for the four point function are simply the bubble chains of Fig.II.5.3. They form a geometric series which can be summed explicitly and replaced by a wavy line. This point of view allows to reorganize the graphs of the expansion as in Fig.II.5.4. Again for “wrong sign”  $g$  the model is asymptotically free. Nevertheless the “bubble chain” model is too simple to keep some interesting features of asymptotic freedom like the loglog behavior analyzed below (unless next to leading order in  $N^{-1}$  is included). This is why in this section we choosed to analyze the less trivial planar series.\*

Although not explicitly soluble like the “bubble chain” model, the planar model remains entirely tractable with perturbative methods because of the following main simplification:

### Theorem II.5.1

There exists some constant  $K$  such that the number of planar Feynman graphs for the  $\phi^4$  theory, counted with their proper multiplicity, is bounded at order  $n$  by  $K^n$  (times a function of the number of their external legs).

The proof of this statement is in [KNN]. An other proof together with an asymptotic analysis of the number of planar diagrams at large order is in [BIPZ]. Here we do not need the precise value of  $K$ .

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\* As an interesting problem, which is open to the author’s knowledge, we suggest to work out these models and their  $N \rightarrow \infty$  limit for tensor fields  $\phi$  with more than 1 or 2 indices (there may be no canonical choice for the vertex), and to search for an analogue of Theorem II.5.1.

Theorem II.5.1 means that the planar restriction trims most of the graphs of ordinary  $\phi^4$  theory and that for instance the bare series with a fixed ultraviolet cutoff has a finite radius of convergence.

Following the remarks at the end of last section, we use as a starting point the perturbation theory with mass renormalization fully performed and with bare coupling constant  $g_\rho$  and bare wave function constant  $a_\rho$ . The Schwinger function  $C_{N,planar}^\rho$  are therefore expressed as:

$$C_{N,planar}^\rho = \sum_{G,\mu \leq \rho} \frac{(g_\rho)^n}{S_{planar}(G)} A_{G,\mu}^{MR} \quad (\text{II.5.3})$$

the sum being performed over planar graphs  $G$ , and  $S_{planar}(G)$  being their combinatoric weight in the planar theory. From now on we forget the subscript *planar* in the rest of this section. The mass-renormalized amplitudes  $A^{MR}$  are defined by:

$$A_{G,\mu}^{MR} \equiv \int \prod_v dx_v \prod_{b \in \mathbf{B}(G)} (1 - \tau_b^{0*}) Z_{G,\mu} \quad (\text{II.5.4})$$

We pass to an effective expansion only for the coupling constant. Hence we derive an analogue of Theorem II.4.1, rather than of the more complicated Theorem II.4.3. (This will be justified only a posteriori when asymptotic freedom will make finite the apparent divergences due to the lack of wave function renormalization). Therefore we really do not need all the apparatus of sect. II.3B to treat bipeds. We can stick to the simpler definitions (II.3.4-5) of external and internal indices. We also do not need the lemmas II.3.1 and II.3.2 for the classification of forests. We simply define  $\mathbf{D}_\mu$  as the forest of all quadrupeds  $q \subseteq G$  which satisfy the almost locality condition  $i_q(\mu) > e_q(\mu)$ . Then we derive the following analogue of Theorem II.4.1:

### Theorem II.5.2

There exist  $\rho + 1$  formal power series  $g_i^\rho$ ,  $i = \rho - 1, \dots, -1$  in  $g_\rho$ , such that (II.5.3) is the same power series in  $g_\rho$  as

$$C_N^\rho = \sum_{G,\mu \leq \rho} \left[ \prod_{v \in G} g_{e_v(\mu)}^\rho \right] \frac{1}{S(G)} A_{G,\mu}^{MR;CCUR} \quad (\text{II.5.5})$$

where

$$A_{G,\mu}^{MR;CCUR} \equiv \int \prod_v dx_v \prod_{b \in \mathbf{B}(G)} (1 - \tau_b^{0*}) \prod_{q \in \mathbf{D}_\mu} (1 - \tau_q^*) Z_{G,\mu} \quad (\text{II.5.6})$$

and

$$g_i^\rho = g_{i+1}^\rho + \sum_{\substack{H \text{ quadruped, } \mu \leq \rho \\ i_H(\emptyset) = i+1}} \frac{1}{S(H)} \left[ \prod_{v \in H} g_{e_v(\mu)}^\rho \right] \dots$$

$$\int \prod_v dx_v \prod_{b \in \mathbf{B}(H)} (1 - \tau_b^{0*}) \prod_{q \in \mathbf{D}_\mu(H), q \neq H} (1 - \tau_q^*) \tau_H Z_{H,\mu} \quad (\text{II.5.7})$$

To establish this formula one follows the same path as for Theorem II.4.1, except that the combinatoric of planar Wick contractions is different from the

ordinary one, so that the combinatoric aspect of the Bogoliubov induction must be checked again. The key point to notice is that inserting a 4 point subgraph at a particular vertex, one must now preserve the cyclic ordering in the plane of the 4 double lines or “ribbons” of the planar vertex [Ril]. Apart from that, the proof is just as before.

To go beyond formal power series, we want to choose first  $g_\rho$  so that the power series  $g_i^\rho$  and (II.5.5) are convergent. They will then define the theory. Using Theorem II.5.1 and the bounds of Sect.II.3, this is relatively easy, provided  $g_\rho$  is very small as  $\rho \rightarrow \infty$ . But if  $g_\rho$  is too small as  $\rho \rightarrow \infty$  we end up on an uninteresting theory with  $g_r = 0$ . Hence the real challenge is to find a clever  $\rho$  dependent ansatz for  $g_\rho$  so that as  $\rho \rightarrow \infty$ , the renormalized coupling  $g_{-1}^\rho \equiv g_r^\rho$  tends to some given small fixed  $g_r$ ; then the limit constructed in this way is not trivial, but corresponds to this prescribed renormalized coupling.

We start with a heuristic search for the right ansatz, and then prove that it actually works.

The smallest possible graph in the recursion relation (II.5.7) is our friend the bubble; the only other connected four point graph with two vertices is  $Q_2$  in Fig.II.5.1, which is 0 after mass renormalization is performed. To second order (II.5.7) reduces therefore to:

$$g_i^\rho \simeq g_{i+1}^\rho + b \sum_{\substack{i_1, i_2 = i+1, \dots, \rho \\ i_{12} \equiv \inf\{i_1, i_2\} = i+1}} [g_{j_{12}}^\rho]^2 \int d^4y C^{i_1}(x, y) C^{i_2}(x, y) \quad (\text{II.5.8})$$

By translation invariance the right hand side is independent of  $x$ ;  $b$  is by definition the planar combinatoric coefficient of the bubble, and we define  $j_{12} = \sup\{i_1, i_2\}$ . For fixed  $i_{12}$  after integration over  $y$  the sum over  $j_{12} \geq i_{12}$  is exponentially decreasing as shown in Sect.II.2; furthermore the mass term in the propagators is also small if  $i \gg 1$ . More precisely it is a simple exercise to check:

**Lemma II.5.1** There exists a numerical constant  $\beta_2$  such that:

$$\beta_2 \log M = b \lim_{i \rightarrow \infty} \lim_{\rho \rightarrow \infty} \sum_{\substack{i_{12} = i+1 \\ j_{12} \leq \rho}} \int d^4y C^{i_1}(x, y) C^{i_2}(x, y) \quad (\text{II.5.9})$$

Furthermore both limits are “exponential”, in the sense that for some small enough  $\delta$ :

$$|\beta_2 \log M - b \sum_{\substack{i_{12} = i+1 \\ j_{12} \leq \rho}} \int d^4y C^{i_1}(x, y) C^{i_2}(x, y)| \leq e^{-\delta(\rho-i)} + e^{-\delta i} \quad (\text{II.5.10})$$

$\beta_2$  may be identified with the second order coefficient of the Callan-Symanzik  $\beta$  function [Ca1][Sy1]; in the planar theory it has the particular value  $\frac{16}{16\pi^2}$ ; this value, smaller than the standard value  $\frac{72}{16\pi^2}$  of  $\beta_2$  for one component  $\phi_4^4$  reflects the fact that there are few contraction schemes which respect planarity and cyclic ordering of the planar vertices.

The sum over  $j_{12}$  being exponentially decreasing it should also be no problem to replace (at least in the second order approximation we are using)  $g_{j_{12}}^2$  by  $g_{i+1}^2$  or even  $g_i g_{i+1}$ . With all these changes the recursion relation (II.5.8) takes the simpler form (we forget superscripts  $\rho$  since  $\rho \rightarrow \infty$  has been taken in (II.5.9)):

$$\frac{1}{g_{i+1}} - \frac{1}{g_i} \simeq \beta_2 \log M \quad (\text{II.5.11})$$

an equation whose exact solution is, in terms of  $g_r = g_{-1}$ :

$$g_i = \frac{g_r}{1 + (i+1)g_r \beta_2 \log M} \quad (\text{II.5.12})$$

Let us assume that the approximate “asymptotically free” behavior  $g_i \sim \frac{\text{const.}}{i}$  deduced from (II.5.12) is correct; it has important consequences.

The first consequence is that only third order terms in the recursion (II.5.7) should be relevant in determining the exact form of the ansatz for  $g_\rho$ . This is because relation (II.5.11) when generalized to higher orders becomes:

$$\begin{aligned} \frac{1}{g_{i+1}} - \frac{1}{g_i} &= \beta_2 \log M + \text{const.} g_{i+1} + O(g_{i+1}^2) \\ &\simeq \beta_2 \log M + \frac{\text{const.}}{i} + O\left(\frac{\log i}{i^2}\right) \end{aligned} \quad (\text{II.5.13})$$

Since  $\frac{1}{i}$  is still a divergent series but  $\frac{\log i}{i^2}$  is not, the correct asymptotic behavior of  $g_i^{-1}$ , starting from a fixed value of  $g_{-1} = g_r$  should be  $(\beta_2 \log M)i + \text{const.} \log i + \text{const.}$ . We call  $\frac{\beta_3}{\beta_2}$  the constant in front of  $\log i$  (indeed  $\beta_3$  turns out to be the third coefficient in the Callan-Symanzik beta function).

The second consequence is that if we perform the effective analysis for the wave function constant, we should derive a recursion relation of type (II.4.25). The leading term would correspond to the second order biped  $B_0$  in Fig.II.5.1 (since mass renormalization kills the first order tadpole), hence to a contribution in  $g_{i+1}^2$ . Therefore we can expect:

$$\delta a_{i+1} - \delta a_i \sim \frac{\text{const.}}{i^2} + O\left(\frac{\log i}{i^3}\right) \quad (\text{II.5.14})$$

and  $\sum_0^\infty \delta a_i$  should be finite (and even small for small  $g_r$ ). This means that a bare ansatz in which  $a_\rho$  is constant in  $\rho$  (e.g. close to 1) is acceptable since it leads to a finite renormalized wave function constant, also close to 1 for small  $g_r$ . We summarize this phenomenon by saying that in an asymptotically free theory of this kind the flow of the wave function constant is bounded (and small for small coupling). This is neither true for the flow of the mass nor for the flow of the coupling constant. An other important aspect of this phenomenon is that the apparent logarithmic divergences associated to dangerous bipeds in  $\mathbf{D}_\mu$  are spurious. Although the second Taylor subtraction corresponding to wave function renormalization has not been performed for these bipeds, their contribution is nevertheless finite since there are at least two vertices for each them, and  $\sum_0^\infty i^{-2}$  converges. But one should be aware that there is some price to pay for that: the index space convergence of the amplitudes in the effective expansion (II.5.5) now requires to

use the decay of the vertex factors  $\prod_v g_{e_v(\mu)}$ , and is no longer exponential in index space, but only power-like when bipeds are present. This requires sometimes additional care, as will happen in the computation of the coefficient  $\beta_3$  below.

The conclusion of this heuristic analysis of asymptotic freedom is to justify, as announced, our use of a formula with no wave function renormalization, and to tell us that in order to land onto a small finite renormalized coupling constant, we should try as an ansatz for  $a_\rho$  a constant  $a$  close to 1 and for  $g_\rho$  the formula:

$$g_\rho = f(\rho, C)^{-1}; \quad f(\rho, C) \equiv (\beta_2 \log M) \rho + \frac{\beta_3}{\beta_2} \log \rho + C \quad (\text{II.5.15})$$

where  $\beta_2$  is defined by (II.5.9),  $\beta_3$  is an other computable constant (defined by (II.5.21-22-24) below) which results from the careful study of the subleading terms in the recursion relation (II.5.7), and  $C$  is a (large) constant, whose value is related to the exact value of  $g_r$  that one wants to obtain. Hence the constants  $a$  and  $C$  play the role of two bare parameters which parametrize the two parameter family of theories one is looking for (remember that from the beginning the third parameter, the mass, is the renormalized one).

To formulate as a precise theorem the construction of the planar theory it is convenient to use complex values of  $a$  and of  $C$  in (II.5.15) and to introduce the half-plane  $H_K = \{C | \text{Re } C > K\}$  and the disks  $D_\delta^\delta = \{g | \text{Re } \frac{1}{g} > \frac{2}{\delta}\}$  and  $D_\delta^1 = \{a; |a - 1| < \delta\}$  of radius  $\delta$  and centers respectively at  $\delta$  and 1 (see Fig.II.5.5). (Of course  $D_\delta^\delta = C_{R/2}$  in the notations of section I.5).

The main result is then summarized by:

### Theorem II.5.3

Let  $g_\rho$  be given by (II.5.15), with  $C \in H_K$  and  $K$  large enough, and  $a_\rho = a$ ,  $a \in D_\delta^1$ , with  $\delta$  small enough. Then the recursive relation (II.5.7) and the effective expansion (II.5.5) are absolutely convergent, uniformly in  $H_K \times D_\delta^1$ . The corresponding sums  $g_i^\rho$  and  $C_N^\rho$  are analytic in  $H_K \times D_\delta^1$  and converge uniformly as  $\rho \rightarrow \infty$  to functions  $g_i$  and  $C_N$ , therefore also analytic in  $H_K \times D_\delta^1$ , which define the theory with no cutoff.

In particular there is a doubly analytic map from  $(C, a)$  to the renormalized coupling and the renormalized wave function constant  $(g_r, a_r)$ ; this map is from  $H_K \times D_\delta^1$  to  $D_{\delta'}^{\delta'} \times D_{\delta'}^1$  for some  $\delta'$ . It can be inverted to a map  $(g_r, a_r) \rightarrow (C, a)$  from a smaller double disk  $D_{\delta''}^{\delta''} \times D_{\delta''}^1$  to  $H_K \times D_\delta^1$  for some  $\delta'' < \delta'$ . In this way the theory can be parametrized by the renormalized parameters  $g_r$ ,  $m_r$  and  $a_r$ ; it is in particular analytic in  $g_r$  for  $g_r$  in  $D_{\delta''}^{\delta''}$  and the functions  $C_N(g_r, m_r, a_r)$  are the Borel sums of the ordinary fully renormalized planar series in  $g_r$  with propagator  $(a_r p^2 + m_r^2)^{-1}$  (see section I.5 for the definition of Borel summability).

We prove this theorem by an induction in index space from  $i = \rho$  to  $i = -1$ . Let us study first in more detail (II.5.7). The only third order graphs for the four point function to be considered are  $Q_3$  and  $Q_4$  in Fig.II.5.1. Hence to third order (II.5.7) leads to:

$$\begin{aligned} g_i^\rho - g_{i+1}^\rho &= b \sum_{i+1=i_{12} \leq j_{12} \leq \rho} (g_{j_{12}}^\rho)^2 \int C^{i_1}(x, y) C^{i_2}(x, y) d^4 y \\ &+ c \sum_{i+1=i_{1234} \leq j_{1234} \leq \rho} g_{j_{12}}^\rho g_{j_{134}}^\rho g_{j_{234}}^\rho \int d^4 y d^4 z \mathbf{R}_{int}^U C^{i_1}(x, y) C^{i_2}(x, z) C^{i_3}(y, z) C^{i_4}(y, z) \end{aligned}$$

$$+d \sum_{i+1=i_{1234} \leq j_{1234} \leq \rho} g_{j_{12}}^\rho g_{j_{1234}}^\rho g_{j_{34}}^\rho \int d^4y d^4z \mathbf{R}_{int}^U C^{i_1}(x, y) C^{i_2}(x, y) C^{i_3}(y, z) C^{i_4}(y, z) \\ + \text{graphs with at least 4 vertices} \quad (\text{II.5.16})$$

where we use obvious notations  $i_{abc\dots} = \inf\{i_a, i_b, i_c, \dots\}$ ,  $j_{abc\dots} = \sup\{i_a, i_b, i_c, \dots\}$ ,  $b$ ,  $c$  and  $d$  are respectively the planar combinatoric coefficients of the graphs  $Q_1$  (the bubble)  $Q_3$  and  $Q_4$  of Fig.II.5.1, and the “useful” internal renormalization operator  $\mathbf{R}_{int}^U$  is by definition  $\prod_{h \in \mathbf{D}_\mu(H), h \neq H} (1 - \tau_h^*)$ . In the second term of (II.5.16) we have therefore:

$$\begin{aligned} & \mathbf{R}_{int}^U C^{i_1}(x, y) C^{i_2}(x, z) C^{i_3}(y, z) C^{i_4}(y, z) \\ &= C^{i_1}(x, y) C^{i_2}(x, z) C^{i_3}(y, z) C^{i_4}(y, z) \text{ if } i_{34} \leq j_{12} \\ & \mathbf{R}_{int}^U C^{i_1}(x, y) C^{i_2}(x, z) C^{i_3}(y, z) C^{i_4}(y, z) = \\ & C^{i_1}(x, y) [C^{i_2}(x, z) - C^{i_2}(x, y) C^{i_3}(y, z) C^{i_4}(y, z)] \text{ if } i_{34} > j_{12} \end{aligned} \quad (\text{II.5.17})$$

In the third term, the full amplitude factorizes as a product so that when one internal operator  $(1 - \tau_h^*)$  is performed, the result vanishes. Therefore:

$$\begin{aligned} & \mathbf{R}_{int}^U C^{i_1}(x, y) C^{i_2}(x, y) C^{i_3}(y, z) C^{i_4}(y, z) \\ &= C^{i_1}(x, y) C^{i_2}(x, y) C^{i_3}(y, z) C^{i_4}(y, z) \text{ if } i_{12} \leq j_{34} \text{ and } i_{34} \leq j_{12} \\ & \mathbf{R}_{int}^U C^{i_1}(x, y) C^{i_2}(x, y) C^{i_3}(y, z) C^{i_4}(y, z) = 0 \text{ otherwise} \end{aligned} \quad (\text{II.5.18})$$

To obtain a recursion relation involving solely  $g_i^\rho$  and  $g_{i+1}^\rho$  at third order, one should replace, in each third order term of (II.5.16) every  $g_k^\rho$  with  $k \geq i+1$  simply by  $g_{i+1}^\rho$ , since they are equal at first order. But in the second order term one should reexpress  $g_{j_{12}}^\rho$  in terms of  $g_{i+1}^\rho$ , taking into account the third order terms that this operation generates. This results in contributions with a main bubble  $H$  corresponding to the initial second order graph (hence with  $i_H = i+1$ ), and a reduction vertex corresponding to the counterterm for an other bubble  $h$ , generated by the second order recursion relation (II.5.8), which we apply to reexpress  $g_{j_{12}}^\rho$  in terms of  $g_{i+1}^\rho$ . This second bubble  $h$  may be inserted “transversally” or “longitudinally”, in which case we associate this contribution respectively to the graphs  $Q_3$  and  $Q_4$ . For  $Q_3$ , the corresponding counterterm completes the renormalization of  $h$ ; namely the inner bubble  $h$  satisfies  $j_{12} \geq i_{34} > i+1 = i_{12}$ : adding its counterterm replaces precisely the operator  $\mathbf{R}_{int}^U$  by  $\mathbf{R}_{int} \equiv (1 - \tau_h^*)$ , the full internal renormalization operator, except for a subtlety; the counterterm for internal subgraphs  $h$  with  $i_{34} = i+1$  is missing. We add it and subtract it, so that  $Q_3$  gets full internal renormalization, and there is an exceptional term which is  $+ \tau_H \tau_h^*$ , with  $i_H = i_h = i+1$ .

Similarly for  $Q_4$ , taking into account the symmetry factor of the graph, missing counterterms for  $h$  to the right and  $j_{12} \geq i_{34} > i+1 = i_{12}$  and for  $h$  to the left and  $j_{34} \geq i_{12} > i+1 = i_{34}$  are generated. They correspond almost exactly to the first case of (II.5.18), so that the contribution for  $Q_4$  becomes 0 as in the second case of (II.5.18), except again for the subtlety that the case  $i_{12} = i_{34}$  remains. But combining this term with the exceptional term for the graph  $Q_3$ , we reconstruct exactly the product of two independent bubble counterterms  $\tau_H \tau_{H'}$  (with their

full combinatoric coefficients), and with  $i_H = i_{H'} = i + 1$ . The conclusion of this tedious analysis is that to third order in  $g_{i+1}$ , (II.5.16) becomes:

$$\begin{aligned}
g_i^\rho - g_{i+1}^\rho &= b \sum_{i+1=i_{12} \leq j_{12} \leq \rho} (g_{i+1}^\rho)^2 \int C^{i_1}(x, y) C^{i_2}(x, y) d^4y \\
&+ c \sum_{i+1=i_{1234} \leq j_{1234} \leq \rho} (g_{i+1}^\rho)^3 \int d^4y d^4z \mathbf{R}_{int} C^{i_1}(x, y) C^{i_2}(x, z) C^{i_3}(y, z) C^{i_4}(y, z) \\
&+ (g_{i+1}^\rho)^3 \left\{ b \sum_{i+1=i_{12} \leq j_{12} \leq \rho} \int C^{i_1}(x, y) C^{i_2}(x, y) d^4y \right\}^2 \\
&+ \text{contributions with at least 4 vertices}
\end{aligned} \tag{II.5.19}$$

It remains to check that the “contributions with at least 4 vertices” in (II.5.19) are unimportant. These corrections correspond to graphs with 4 vertices or more in (II.5.16) or are generated by the recursion which changes (II.5.16) into (II.5.19). Taking into account the form of the error terms in (II.5.13) and (II.5.10), we expect a bound of the type  $O(\frac{\log i}{i^4}) + i^{-2}O(e^{-\delta(\rho-i)})$  for the sum of all these terms. Such a bound would not be too hard to prove if one could directly combine the uniform bound on usefully renormalized graphs of Section II.3 with Theorem II.5.1. Indeed taking  $C$  large enough, one is certainly inside the convergence radius of the planar power series; the remainder in (II.5.16) for  $i+1=\rho$  can be evaluated by  $O(g_\rho^4) = O(\rho^{-4})$ . By induction, one would remain inside convergence radius of the series in (II.5.16) for smaller  $i$ 's; the behavior (II.5.13) would be checked inductively and the series of contributions generated from (II.5.16) to (II.5.19) again would be controlled inductively. For this last bound one can simply use the exponential decay in index space generated by phase space analysis to obtain a uniform bound  $O((g_{i+1}^\rho)^4) = O(i^{-4})$  for series of contributions of the type  $\sum_{j>i+1} e^{-\delta(j-i)} \sum_{n \geq 4} (g_j^\rho)^n a_n$ .

This program is basically right except for one subtle point: the amplitudes appearing in (II.5.16) are not exactly the usefully renormalized ones; in particular no wave function renormalization is performed. We need to use the decay of the effective constants to sum over logarithmically divergent mass-renormalized bipeds. This is possible because there are at least two effective constants available for each such biped. Hence the analysis sketched above remains correct, except that graphs  $H$  with  $k$  bipeds,  $n$  vertices and  $i_H > i$  should not be evaluated naively by the regular bound  $O(i^{-n})$  but rather by  $O(i^{-(n-k)})$  since each logarithmic divergence “eats” the decay of one coupling constant. This small change would not be relevant at all, except for the fact that the smallest possible  $H$  with  $k \neq 0$ , the graph  $Q_5$  in Fig.II.5.1, should not be considered as fourth order, but rather promoted to third order, so that it becomes relevant for the correct value of  $\beta_3$  \*.

The detailed analysis of the graph  $Q_5$  is easy. The non trivial biped  $B_0$  corresponds to lines 4,5,6. The logarithmic divergence is fully contained in the case  $i_{456} > i_{123}$  and the corresponding (mass-renormalized) contribution is obtained by

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\* Of course the phenomenon discussed here at the level of graphs is well known at the level of the Callan-Symanzik equation; the second order term in (II.5.14) reacts on the third order term of this equation.

applying  $\tau_{B_0}^{1*}$ . Hence if  $e$  is the combinatoric coefficient for the planar graph  $Q_5$  the full apparent third order contribution of  $Q_5$  in (II.5.16), up to  $O(\frac{\log i}{i^2})$ , is simply:

$$\begin{aligned} & e (g_{i+1}^\rho)^2 \sum_{i_{123}=i+1} \sum_{k=i_{456}>i_{123}} \int d^4x d^4y d^4z C^{i_1}(u, z) \frac{|y-x|^2}{8} \\ & \{C^{i_2}(x, u) \Delta C^{i_3}(x, z)\} \frac{1}{(\beta_2 k \log M)^2} C^{i_4}(x, y) C^{i_5}(x, y) C^{i_6}(x, y) \end{aligned} \quad (\text{II.5.20})$$

We generalize now Lemma II.5.1 to the contributions in (II.5.19-20):

**Lemma II.5.2** There exists numerical constants  $\gamma_3$  and  $\delta_3$  such that:

$$\begin{aligned} \gamma_3 \log M &= c \lim_{i \rightarrow \infty} \lim_{\rho \rightarrow \infty} \sum_{i+1=i_{1234} \leq j_{1234} \leq \rho} \int d^4y d^4z \dots \\ & \mathbf{R}_{int} C^{i_1}(x, y) C^{i_2}(x, z) C^{i_3}(y, z) C^{i_4}(y, z) \end{aligned} \quad (\text{II.5.21})$$

$$\begin{aligned} \frac{\delta_3}{\beta_2} &= e \lim_{i \rightarrow \infty} i \cdot \lim_{\rho \rightarrow \infty} \sum_{i_{123}=i+1} \sum_{k=i_{456}>i+1} \int d^4x d^4y d^4z \dots \\ & C^{i_1}(u, z) \frac{|y-x|^2}{8} \{C^{i_2}(x, u) \Delta C^{i_3}(x, z)\} \frac{1}{(\beta_2 k \log M)^2} C^{i_4}(x, y) C^{i_5}(x, y) C^{i_6}(x, y) \end{aligned} \quad (\text{II.5.22})$$

Furthermore the sum in (II.5.21) converges exponentially and the generalization of (II.5.10) holds, namely differences between the right hand side of (II.5.21) at finite  $i$  and  $\rho$  and the left hand side are uniformly bounded by  $O(e^{-\delta(\rho-i)}) + e^{-\delta i}$ .

The limit in (II.5.22) is more complex; up to exponentially small errors in  $\rho - i$  and  $i$ , we have an integral invariant under translation in index space, hence a contribution  $\text{const.} \sum_{i+1 < k \leq \rho} k^{-2} = \text{const.}[i^{-1} - \rho^{-1} + O(i^{-2})]$ .

Finally by Lemma II.5.1 the third contribution in (II.5.19) converges simply to  $(\beta_2 \log M)^2$  with exponentially small corrections as  $\rho$  and  $i$  tend to infinity. Therefore we obtain:

**Lemma II.5.3**

$$\begin{aligned} g_i^\rho - g_{i+1}^\rho &= \beta_2 \log M (g_{i+1}^\rho)^2 + \gamma_3 \log M (g_{i+1}^\rho)^3 + (\beta_2 \log M)^2 (g_{i+1}^\rho)^3 \\ & + \frac{\delta_3}{\beta_2} (g_{i+1}^\rho)^2 \left( \frac{1}{i} - \frac{1}{\rho} \right) + O\left(\frac{\log i}{i^4}\right) + (g_{i+1}^\rho)^2 O(e^{-\delta(\rho-i)}) \end{aligned} \quad (\text{II.5.23})$$

Let us define:

$$\beta_3 = \gamma_3 + \delta_3 \quad (\text{II.5.24})$$

The behavior of the recursion relation (II.5.23) is investigated easily under the initial ansatz (II.5.15), even for complex  $C$  with  $\text{Re}C$  large.

**Lemma II.5.4** For  $\text{Re}C$  large enough, we have:

$$|\frac{1}{g_i^\rho} - (\beta_2 \log M)i + \frac{\beta_3}{\beta_2} \log i + C| \leq \text{Re}C \quad (\text{II.5.25})$$

This is because we may rewrite (II.5.23) as:

$$\frac{1}{g_{i+1}^\rho} - \frac{1}{g_i^\rho} = \beta_2 \log M + \frac{\beta_3}{\beta_2 i} - \frac{\delta_3}{\beta_2 \rho} + O\left(\frac{\log i}{i^2}\right) + O(e^{-\delta(\rho-i)}) \quad (\text{II.5.26})$$

(II.5.25) follows from the uniform summability of  $\frac{\log i}{i^2}$ ,  $e^{-\delta(\rho-i)}$ , and the obvious bound:

$$\sum_{j=i}^{\rho} \frac{\delta_3}{\beta_2 \rho} \leq \frac{\delta_3}{\beta_2} \quad \forall i \quad (\text{II.5.27})$$

We can now extend uniform bounds like Theorem II.3 to the amplitudes which appear in the particular expansion II.5.3:

**Lemma II.5.5** There exists a constant  $K$  such that:

$$|\sum_{\mu} [\prod_{v \in G} g_{e_v(\mu)}^\rho] A_{G,\mu}^{MR;CCUR}| \leq K^{n(G)} (1 + \sup_j |p_j|)^{\hat{N}} \quad (\text{II.5.28})$$

This bound is obtained by combining the bounds on usefully renormalized amplitudes (Theorem II.3) for the pieces of  $G$  without bipeds, the argument that mass insertions do not create renormalon effects (see Theorem II.3.2) and the existence of at least two specific coupling constants of the right scale associated to each dangerous biped; using the decay of these constants as expressed in Lemma II.5.4 the logarithmic divergence for these bipeds not only becomes convergent but does not disturb the uniform nature of the estimate.

It is important to notice that the proof of Lemmas II.5.3-4-5 is inductive. For each scale  $i$  from  $\rho$  to  $-1$ , Lemma II.5.3 is proved first, then Lemma II.5.4, then the piece of Lemma II.5.5 which deals with the subgraphs  $G_i^k$  of  $G$  (made of lines with indices  $j \geq i$ ). There is no logical loop, because the uniform bounds on the remainders necessary for Lemma II.5.3 at scale  $i$  only depend on the bounds (II.5.26) for effective couplings of scales  $j > i$ , and on bounds of the type (II.5.28) for subgraphs  $H$  with  $i_H > i$ .

Once the renormalization group “discrete flow” for the coupling constants  $g_i$  and the bounds (II.5.28) are established by Lemmas II.5.4-5 for all scales, it is easy to complete the proof of Theorem II.5.3. Using Lemmas II.5.4, II.5.5 and Theorem II.5.1, the series (II.5.5) are absolutely and uniformly convergent, hence their sum defines the planar theory with cutoff  $\rho$ . Furthermore every estimate and the sums over index space being uniform in  $\rho$ , the effective constants  $g_i^\rho$  have limits  $g_i$  as  $\rho \rightarrow \infty$  which still satisfy Lemmas II.5.4-5. This constructs the theory without cutoff. The analyticity result follows in the straightforward way from uniform convergence of series term by term analytic. The results on the behavior of  $g_r$  and  $a_r$  as functions of  $C$  and  $a$  are obtained by considering the recursion relation (II.5.14) for the wave function constant in addition to the flow of the coupling constant expressed in Lemma II.5.3.

Finally for Borel summability in  $g_r$ , one checks directly the hypotheses of Nevanlinna-Sokal theorem (Theorem I.5.1). The region of Fig.II.5.5 is exactly the region necessary to apply this theorem (the true region of analyticity is in fact much larger because the discrete flow (II.5.23) still leads to an asymptotically free theory ( $g_\rho \rightarrow 0$  as  $\rho \rightarrow \infty$ ) for  $g_r$  at least in the region pictured in Fig.II.5.6) The

uniform Taylor remainder estimates (I.5.2) are just an other exercise in establishing factorial bounds (in this case solely due to the renormalon effects [Ri1]). Since this kind of bounds has been studied at length in Sect.II.3, we leave this problem to the reader.

It might be interesting at this stage to compare briefly what we have done to the standard continuous renormalization group flows, limiting ourselves for simplicity to the example of the  $\beta$  function. In the standard definition of the Callan Symanzik or renormalization group equations the key rôle is played by the ultraviolet  $\beta_{uv}$  function which is defined as  $\frac{dg_{ren}}{dx}|_{g_\rho}(g_{ren})$ , where  $g_{ren}$  is the renormalized coupling;  $g_\rho$ , the bare coupling, is held fixed, and  $\frac{x=\log \kappa}{m_{ren}}$  is the logarithm of the quotient between the scale  $\kappa$  of the ultraviolet cutoff and the renormalized mass, hence in our case where we work with a unit renormalized mass which is fixed,  $x$  is equal to  $\rho \log M$ . In the BPHZ scheme that we use, it is known that there are formulas which relate this ultraviolet  $\beta$  function to renormalized Schwinger functions  $\Gamma_\Delta^{(N)}$  at zero momentum with one mass insertion  $\Delta$  on one propagator and the minimal subtraction prescription. This means that two or four point functions which contain the mass insertion  $\Delta$  should be less subtracted, according to their true (improved) degree of convergence. The formula for the ultraviolet  $\beta$  formal power series is [IZ]:

$$\beta_{uv}(g) = \left[ \frac{1}{1 + \frac{\partial \Gamma_\Delta^2(0)}{\partial p^2}} \right] [-2\Gamma_\Delta^{(4)}(0) + 4g \frac{d\Gamma_\Delta^2(0)}{dp^2}] \quad (\text{II.5.29})$$

This formula can be used for practical numerical computations of coefficients such as  $\beta_2$  or  $\beta_3$ . What is its relationship to the method described above? When we compute the difference  $g_i^\rho - g_{i+1}^\rho$  we compute clearly a discrete analogue of the ultraviolet beta function with cutoff  $\kappa = M^\rho$ , but this function is not expressed as a power series in the last (renormalized) coupling  $g_i^\rho$  but in the whole sequence of previous effective couplings. Usually this is a better way of doing than expressing it in terms of a single renormalized constant, an operation which generates useless counterterms and renormalons. However in this particular case, to study the behavior of the effective couplings it is practical to develop at least the first orders of the equation in terms of the last coupling  $g_i^\rho$ , as is done above to third order. In this way we see that the coefficients of the usual  $\beta$  series (II.5.29) are not exactly generated both because of the remaining ultraviolet cutoff and of a slice effect due to the fact that our flow is discrete rather than continuous. More precisely the condition that one propagator in our contributions to the flow is in slice  $i$  is asymptotically the analogue of the mass insertion in (II.5.29). However since our slices have finite thickness, there are some terms with several legs in the slice (a situation of measure zero for infinitesimal slices hence for continuous flows). These terms are the source of corrections to the  $\beta$  function such as the term in  $(\beta_2 \log M)^2$  in (II.5.23) which are characterized by a power of  $\log M$  higher than one. In other words if we develop the finite difference equation (II.5.23) up to a given finite order in terms of the last (renormalized) constant  $g_i^\rho$ , take  $\rho \rightarrow \infty$ , divide by  $\log M$  (since  $x = \rho \log M$  and take  $M \rightarrow 1$ , then order by order only the regular contribution of the ultraviolet  $\beta$  function will survive, but in a scheme with an infrared cutoff; in the limit  $i \rightarrow \infty$  the small corrections to scale invariance due to the mass disappear and we will find exactly the same coefficients than in (II.5.29) (this explains the two limits in (II.5.9)).

In conclusion the discrete flows considered in this section and later in part III are naturally expressed in terms of effective quantities with effective constants; they are the correct way to replace the formal renormalization group functions by well defined ones. The standard formal power series are recovered when cutoffs are removed and effects due to the discretization are removed.

## II.6. The large order behavior of perturbation theory

*Aussi loin que la science recule ses frontières,  
 et sur tout l'arc étendu de ces frontières,  
 on entendra courir encore la meute chasseresse du poète.*  
 – Saint-John-Parseval

In this last section on perturbation theory we no longer discuss large order bounds for individual Feynman amplitudes, but consider the more difficult problem of the exact large order *behavior* of the renormalized perturbation series. We will meet again the problem of the large number of graphs in the ordinary  $\phi^4$  theory, and the renormalon problem for  $\phi_4^4$ , and discuss how they shape the large order behavior of the theory, using the convenient mathematical formalism of the Borel transform introduced in section I.5. The rigorous results obtained so far are still fragmentary and in our opinion a lot of interesting work remains to be done in this area.

In the regular  $\phi^4$  theory the total number of Wick contractions for graphs (not necessarily connected) with  $n$  vertices and a fixed number  $N$  of external lines is  $(4n + N - 1)!!$ . The number of connected graphs is of course smaller, but it is rather easy to show that it is more than  $(\text{const.})^n n!^2$ ; indeed we may first build a spanning tree in more than  $(\text{const.})^n n!$  different ways (apply Cayley's theorem, Sect.I.4, with coordination numbers bounded by 4), then still have  $(\text{const.})^n n!$  different contraction schemes for the remaining lines.

In a theory like  $\phi_1^4$  (the anharmonic oscillator), all graphs add up at any given order with the same sign. Furthermore it is easy to check that any amplitude of order  $n$  satisfies both an upper and a lower bound of the kind  $(\text{const.})^n$ . The upper bound is simply Weinberg's uniform theorem, although in this case it may be obtained for instance in  $\alpha$  space by simpler arguments. The lower bound is also extremely easy in  $\alpha$  space: simply integrate only over  $1 \leq \alpha_i \leq 2 \forall i$  and use the fact that the number of spanning trees is bounded by  $2^{l(G)} \leq 4^{n(G)}$  for an upper bound on the Symanzik polynomial  $U_G$ .

Hence in this case of  $\phi_1^4$ , taking into account the factor  $\frac{1}{n!}$  in (I.4.1), the  $n$ -th order of perturbation theory  $a_n$  satisfies some bound of the type:

$$K_1^n n! \leq (-1)^n a_n \leq K_2^n n! \quad (\text{II.6.1})$$

for some positive constants  $K_1$  and  $K_2$ . As a consequence, the radius of convergence of the perturbative series is 0; we say in short that it diverges. Nevertheless

a behavior like (II.6.1), although incompatible with ordinary summability may still allow Borel summability.

It is not easy to prove that the corresponding renormalized series diverge for  $\phi^4$  in higher dimensions. BPHZ renormalization in two dimensions is equivalent to Wick ordering, hence simply suppresses the graphs with tadpoles. It has been first proved in [Ja] that enough graphs remain so that the n-th order of perturbation series for  $\phi_2^4$  still satisfy a lower bound which implies divergence. For  $\phi_3^4$  there is a first non trivial mass renormalization which changes the sign of some amplitudes. Lower bounds for sums with different signs is more difficult because one has to rule out the possibility of systematic cancellations. The divergence for  $\phi_3^4$  was proved in [dCR2], but no longer from a lower bound simply on the n-th order of the series; the argument already mixes different orders. For  $\phi_4^4$  a rigorous proof is still missing; renormalization introduces changes of signs much more difficult to track, so that a proof of divergence seems impossible except as a by-product of a detailed analysis of the large order behavior. Significant progress towards this goal has been made in [MNRS], [DFR], using the multiscale representation, and the main goal of this section is to introduce the reader to this approach.

We turn now to a description of the rigorous results and heuristic expectations on the large order behavior of  $\phi^4$ . We will not comment on the large order behavior of other models, except to notice that theories with fermionic fields have better convergence properties due to the cancellations in the corresponding fermionic determinants; this fact, a direct consequence of Pauli's principle, will be used extensively in the construction of the Gross-Neveu model (Sect.III.4). The study of the corresponding large order *behavior* of these models is much less advanced than for bosonic theories, and is a beautiful open problem.

The large order behavior of  $\phi^4$  is of course best understood in the one dimensional case (anharmonic oscillator) (not to speak of the 0-dimensional case, which is the study of the moments of the measure  $e^{-gx^4-\alpha x^2} dx$ , for which explicit formulae can be derived in terms of hypergeometric functions [Wig2]). There is a long history of numerical and rigorous results on the anharmonic oscillator, using BKW methods or the functional integral and steepest descent methods, reviewed in [Si2]. But here we intend to put the emphasis on higher dimensions, d=3 and mostly d=4. The first major progress in these cases came from a semi-rigorous extension of the steepest descent to the functional integral of  $\phi_d^4$ , called the Lipatov method [Lip]. This method was developed and applied in [BGZ]. Let us summarize now its guiding principle.

For simplicity the large order behavior may be investigated on the simplest typical quantity for the theory with non trivial renormalized perturbation theory. For instance in 1,2 and 3 dimensions we may use the pressure:

$$p = \lim_{\Lambda \rightarrow \infty, \rho \rightarrow \infty} \frac{1}{|\Lambda|} \log Z_{\rho, \Lambda} = \sum (-g)^n a_n \quad (\text{II.6.2})$$

but in 4 dimensions the renormalized series for the pressure or for the 2 and 4 point functions at 0 external momenta are trivial, so that the simplest quantity may be the connected 6 point function at 0 external momenta, or the 2 point function at a particular momentum, or the connected 4 point function at some symmetric set

of external momenta. If  $C_N$  is the quantity of interest we write again:

$$C_N = \sum_n (-g_R)^n a_n^R \quad (\text{II.6.3})$$

with  $g_R$  the renormalized constant.

The Lipatov method concludes to an asymptotic behavior of  $a_n$  or  $a_n^R$  at large  $n$  which is always of the type:

$$a_n \simeq n! a^n \cdot n^b \cdot c(1 + O(1/n)) \quad (\text{II.6.4})$$

where  $a$ ,  $b$  and  $c$  are some constants.  $a$  depends only on the dimension (and may also depend on some parameters of the theory, like the mass for  $d < 4$ ).  $b$  may also depend on which particular Schwinger function one is investigating (hence, may depend on  $N$ ) and  $c$  depends on further details of the theory (in particular in 4 dimensions, on the particular renormalization scheme and subtraction scale which is used).

Since we are interested in universal features of the large order behavior which are valid for the perturbative series of any reasonable quantity in the theory,  $a$  is the most important constant in (II.6.3). The Lipatov method predicts:

$$a = e^{-\inf_{\phi} S(\phi)+2} \quad (\text{II.6.5})$$

$$S(\phi) = \frac{1}{2} \int (\partial_\mu \phi \partial^\mu \phi + m^2 \phi^2) - \log \int \phi^4 \quad (\text{II.6.6})$$

where the infimum in (II.6.5) is taken over the appropriate Sobolev space where the functional (II.6.6) is well defined.

Before providing a heuristic motivation for (II.6.5-6), let us show that the functional  $S(\phi)$  is bounded below in dimensions  $d \leq 4$  by virtue of some Sobolev inequality. In dimension  $d \leq 4$  indeed there exists a constant  $K_1$  (depending on  $m$  for  $d < 4$ ) such that:

$$\int_{\mathbb{R}^d} \phi^4 \leq [K_1 \int_{\mathbb{R}^d} (\partial_\mu \phi \partial^\mu \phi + m^2 \phi^2)]^2 \quad (\text{II.6.7})$$

so that whenever  $\phi$  belongs to the Sobolev space  $H^{1,2}$  in which the right hand side of (II.6.7) is well defined, it belongs also to  $L^4$ . We call  $K_1(d)$  the infimum of the constants  $K_1$  for which (II.6.7) holds; this infimum is of course a minimum, i.e. (II.6.7) still holds with  $K_1(d)$  instead of  $K_1$ .

Similarly for  $C^\infty$  functions with compact support (or under suitable decay conditions at infinity) we have also for some optimal  $K_2(d)$ :

$$\int_{\mathbb{R}^d} \phi^4 \leq [K_2(d) \int_{\mathbb{R}^d} (\partial_\mu \phi \partial^\mu \phi)]^2 . \quad (\text{II.6.8})$$

Remark that  $a$  in (II.6.5) and  $K_1(d)$  are related by:

$$a = (4K_1(d))^2 \quad (\text{II.6.9})$$

since taking  $\phi = \alpha f$  with  $\int (\partial_\mu f \partial^\mu f + m^2 f^2) = 1$  we may optimize in (II.6.6) and get:

$$\inf_{\alpha, f} \left( \frac{\alpha^2}{2} - \log \alpha^4 \int f^4 \right) = \inf_f (2 - \log 4^2 \int f^4) = 2 - \log(4K_1(d))^2 \quad (\text{II.6.10})$$

In dimension 1, 2 and 3 the smallest constants  $K_1(d)$  and  $K_2(d)$  for which (II.6.7) or (II.6.8) hold are different;  $K_1(d)$  depends on  $m$ , and for  $m \neq 0$ ,  $K_1(d) < K_2(d)$ . But in the critical case  $d = 4$  we have equality:  $K_1(d) = K_2(d)$ . In fact in dimensions 1, 2 and 3 the infimum in (II.6.5) is a minimum, hence is attained for a particular smooth minimizing function  $\phi_0$  which by the variational principle is a solution of the differential equation:

$$(-\Delta + m^2)\phi_0 = \lambda \frac{\phi_0^3}{\int \phi_0^4} \quad (\text{II.6.11})$$

for some constant  $\lambda$ ; we may find a particular  $\phi_0$  with radial symmetry (to break the translation invariance of (II.6.11)) which has fast decrease at infinity, and  $S(\phi_0)$  depends on  $m$ . This is also true in 4 dimensions for the finite volume analogue of  $S(\phi)$ ,  $S_{\Lambda, X}(\phi)$  (one has to specify some boundary conditions  $X$  on  $\partial\Lambda$ ), which attains its minimum for a particular smooth  $\phi_{0,\Lambda,X}$ , and the corresponding minimum  $K_{1,\Lambda,X}(d=4)$  of  $S_\Lambda$  depends on  $\Lambda$  and  $m$ . As  $\Lambda \rightarrow \infty$ ,  $K_{1,\Lambda,X} \rightarrow K_1(d=4)$ ; it can be shown however that the minimizing functions do not converge, that  $K_1(d=4) = K_2(d=4)$ , the optimal constant for which (II.6.8) holds, which is of course independent of  $m$ , and that the infimum of  $S(\phi)$  is therefore a true infimum (not attained). Furthermore the value of this infimum may be computed exactly:  $K_1(d=4) = K_2(d=4) = \frac{1}{4\pi} \sqrt{\frac{3}{2}}$  [Aub], hence  $a = \frac{3}{2\pi^2}$  in (II.6.5). This change of behavior in dimension 4 is of course due to the marginal character of the Sobolev inequality (II.6.7) in this case.

A crude motivation for the Lipatov prediction of the value (II.6.5) of  $a$  in (II.6.4) goes as follows. Let us pretend that an ultra violet cutoff  $\kappa_n$  and a finite volume  $\Lambda_n$  (with boundary conditions  $X_n$ ) may be imposed on the  $n$ -th order of perturbation theory in such a way that  $\kappa_n, \Lambda_n \rightarrow \infty$  as  $n \rightarrow \infty$ , but that as long as the leading order behavior is concerned,  $a_n^R$  and  $a_n^{X_n, \Lambda_n, \kappa_n}$  (the bare amplitudes with cutoffs and *no* renormalization) are equivalent. (We will see below that this assumption turns out to be justified in dimensions less than 4 but wrong for  $d=4$ ).

Let us assume also that the restriction of connectedness in  $a_n$  is also unimportant as far as leading large order behavior is concerned. Then the leading behavior of  $(-1)^n a_n^R / n!$  is the same as the one of

$$b_n = \int \left( \frac{1}{n!} \right)^2 \left( \int_{\Lambda_n} \phi^4 \right)^n d\mu_{X_n, \Lambda_n, \kappa_n} \quad (\text{II.6.12})$$

For conceptual simplicity, let us choose the ultraviolet cutoff to be a lattice with lattice spacing  $\delta_n$  such that  $\lim_{n \rightarrow \infty} \delta_n = 0$ , and let us choose the dependence in  $n$  of  $\Lambda_n$  so that the total number of sites in  $\Lambda_n$  grows only slowly, like  $n^\epsilon$ . The gaussian measure on the lattice may be written in terms of the ordinary Lebesgue measure (see (I.3.13)). Rescaling  $\phi$  to  $\psi = \frac{\phi}{\sqrt{n}}$  we get:

$$\begin{aligned}
b_n &= \left(\frac{n^n}{n!}\right)^2 \int e^{-(n/2)\{(\partial_\mu \phi \partial^\mu \phi)/n) + m^2(\phi^2/n) - 2 \log \int_{\Lambda_n} (\phi^4/n^2)\}} \prod_{x \in \Lambda_n} d\phi(x) \\
&= (\sqrt{n})^{n^\epsilon} \left(\frac{n^n}{n!}\right)^2 \int e^{-(n/2)\{(\partial_\mu \psi \partial^\mu \psi) + m^2 \psi^2 - 2 \log \int_{\Lambda_n} \psi^4\}} \prod_{x \in \Lambda_n} d\psi(x)
\end{aligned} \quad (\text{II.6.13})$$

so that:

$$(b_n)^{1/n} = e^{n^\epsilon \frac{\log n}{2n}} \left(\frac{n^n}{n!}\right)^{2/n} \|e^{-S(\phi)}\|_n \quad (\text{II.6.14})$$

which, as  $n \rightarrow \infty$  should tend to  $1.e^2 \|e^{-S(\phi)}\|_\infty$ , hence to  $a$  by (II.6.5).

To refine this crude prediction, one has to apply the steepest descent method to the functional integral (II.6.12), hence to expand the action around the configurations which minimize  $S(\phi)$ . By a rescaling of  $\phi$  the variational equation (II.6.11) is the same as the equation of motion of the theory  $(-\Delta + m^2)\phi = g\phi^3$ , and these configurations are the non trivial classical solutions of finite action which are called instantons. Expanding around these solutions to second order, as usual, and performing the corresponding gaussian functional integral gives an explicit determinant with which one can compute subleading coefficients at large order like  $b$  and  $c$  in (II.6.4), and in principle a systematic expansion in  $1/n$ . The Lipatov method in this way analyzes the asymptotic behavior of perturbation theory at large order and relates it to singularities on the negative real axis in the Borel plane which are therefore also called *instanton* singularities. In the case of the anharmonic oscillator these singularities are known to induce a corresponding essential singularity at  $g = 0$  corresponding to a cut along the negative real axis for quantities like the ground state energy (II.6.2). We do not develop further this point of view here, refering to [Zin] for a review on instanton calculus.

However in the standard Lipatov argument as sketched above, the possible effect of renormalization on large order behavior is neglected; more precisely it is assumed simply to change the determinants corresponding to fluctuation around the saddle points into renormalized determinants. This assumption is expected to fail in the case of  $\phi_4^4$ , where renormalization affects in a major way the large order behavior. As argued by Parisi and 'tHooft [Pa1-2][tH3], the factorial behavior of single Feynman graphs that we met and discussed at length in Section II.3 creates corresponding singularities on the right hand side of the real axis in the Borel plane, called *renormalons*. It happens for ordinary  $\phi^4$  as well as for vector  $\phi^4$  models with  $N$  components that the first renormalon singularity on the positive real axis is closer to the origin of the Borel plane than the first instanton singularity on the negative real axis, so that the large order behavior of  $\phi_4^4$  is in fact expected to be governed more by renormalization than by the instantons of the Lipatov method. For instance for one component  $\phi_4^4$  the position of the first expected renormalon is at  $t = 2/\beta_2 = \frac{4\pi^2}{9}$ ; in contrast the value  $a_{Lip} = \frac{3}{2\pi^2}$  above for the Lipatov behavior corresponds to an instanton singularity in the Borel plane at  $t = -\frac{2\pi^2}{3}$  hence farther from the origin by a factor  $3/2$ .\*

The existence of renormalons, if confirmed rigorously, would mean that the  $\phi_4^4$  series are not Borel summable, in contrast to what has been proved for  $\phi_{1,2,3}^4$

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\* There is no deep explanation yet for this simple rational factor of  $3/2$ .

[GGS],[EMS],[MS]. This fact is presumably related to the difficulty in defining a non trivial  $\phi_4^4$  theory satisfying the axioms [Aiz][Frö]. We return to this point in Sect.III.4 with a weak coupling triviality theorem; for reviews on triviality, see [Sok2][GaRi].

We recall now briefly the rigorously proved large order results. For  $\phi_1^4$ , the anharmonic oscillator, Borel summability [GGS] and all of formula (II.6.4) is proven (with expected values of  $a$ ,  $b$  and  $c$ ) [HS]; the Lipatov method has been justified [Sp3][Bre], and much is known about the analyticity properties of the corresponding sums in the Borel plane, either numerically or rigorously (see [Si2],[Wig2]). In two and three dimensions Borel summability has been proved by constructive theory [EMS][MS] and the Lipatov method has been justified basically to leading order only, i.e. up to the computation of the coefficient  $a$  in (II.6.4) [Bre][MR][FR]. Finally in 4 dimensions the results [MNRS][DFR] fall short from proving that the large order behavior is really governed by the first renormalon. Many of these rigorous results are therefore summarized by:

### Theorem II.6.1

- a) The perturbation series for  $\phi_{1,2,3}^4$  are Borel summable (in the Nevanlinna-Sokal sense or in Watson's sense [GGS][EMS][MS]).
- b) For  $\phi_{1,2,3}^4$  there is a disk of analyticity in the Borel plane of radius  $a^{-1}$ , where  $a$  is defined by (II.6.5) and there is a singularity in the Borel plane at  $t = -a^{-1}$  ("instanton") [Bre][MR][FR].
- c) For  $\phi_4^4$  there is a disk of analyticity in the Borel plane of radius  $\frac{2}{\beta_2} = \frac{4\pi^2}{9}$  (the optimal expected disk)[MNRS][DFR].

In dimensions 3 and 4 the comparison between expected and proven results is sketched in Fig.II.6.1.

The proof of item a) relies on constructive methods and we will return to it in the next chapter; more precisely we will give a construction and prove a Borel summability result for infrared  $\phi_4^4$  which is general enough to apply directly with straightforward modifications to ultraviolet  $\phi^4$  in lower dimensions. Such a method is certainly a bit of an overkill for the cases of dimensions 1 and 2, perhaps even for dimension 3, but it is certainly not more complicated than the sum of the specific proofs derived earlier for these lower dimensions, which did not use the full machinery necessary in dimension 4 (multiscale expansion).

Similarly we will not discuss the proof of item b) which concerns superrenormalizable theories because the phase space language is not strictly speaking necessary there and because we prefer to concentrate on the more difficult issue of the four dimensional case. Let us simply state that item b) may be decomposed into the proof of an upper and a lower bound of the Lipatov type. The upper bound [Bre][MR] may be considered a simple corollary of the four dimensional upper Lipatov bound which is discussed below (in dimensions 1,2 and 3, the usefully renormalized series which appear below, in which mass renormalization may be fully performed as in Theorem II.3.2, coincide with the regular series since there are only a few mass renormalizations). For the lower bound, which implies the existence of a singularity at  $t = -a^{-1}$  we refer to [Bre][MR][FR] and notice simply that in three dimensions the changes of signs induced by the single nontrivial mass renormalization require a separate argument in which all orders of perturbation

theory are mixed [FR]. In this sense the proven existence of an “instanton” singularity is still a weaker result than strict asymptotics of the Lipatov type for  $a_n$ , which remains an open problem in dimension 3, even at leading order. Subleading behavior (with the right constants  $b$  and  $c$ ) remains also an open problem, even for  $\phi_4^4$ .

We turn now our attention to item c) of Theorem II.6.1 which we rewrite in more detail as:

### Theorem II.6.2

There exists a function  $\epsilon(n)$  which tends to 0 as  $n \rightarrow \infty$  such that:

$$|a_n^R| \leq n! \left[ \frac{\beta_2}{2} \right]^n (1 + \epsilon(n))^n \quad (\text{II.6.15})$$

where  $a_n^R$  is the n-th order of perturbation theory for the  $\phi_4^4$  model, and  $\beta_2 = 9/2\pi^2$  is the one loop coefficient of the  $\beta$  function (see Sect.II.4-5).

The proof of course generalizes to the  $N$ -component  $\phi_4^4$  theory, in which case  $\beta_2 = (N+8)/2\pi^2$ .

Our goal is to describe in some detail the proof of this theorem. We will skip many of the technicalities, for which the reader is referred to the original articles, but we will try to explain clearly the structure of every important argument. We start by the following bound, proved in [MNRS]:

### Theorem II.6.3: Upper Lipatov bound

There exists a function  $\epsilon(n)$  which tends to 0 as  $n \rightarrow \infty$  such that:

$$|a_n^{UR}| \leq n! [a_{Lip}]^n (1 + \epsilon(n))^n \quad (\text{II.6.16})$$

and  $a_n^{UR}$  is the sum of all usefully renormalized amplitudes (the theorem applies also to  $a_n^{MR,UR}$ , the sum of the mass-renormalized, usefully renormalized amplitudes of Theorem II.3.2).

Upper bounds of the Lipatov type are indeed very natural because of one key simple observation: the critical fields which minimize the functional (II.6.6) also saturate the corresponding Sobolev inequality (II.6.7) (see (II.6.10)). Therefore using the Sobolev inequality on the vertices  $\int \phi^4$  in (II.6.12) should lead to an upper bound of the Lipatov type with correct value of  $a$ . In contrast a lower bound of the Lipatov type typically requires a more complicated analysis on the speed at which, at large order, functional integrals like (II.6.13) become peaked around these minimizing configurations.

However using a Sobolev inequality in (II.6.12) replaces a  $\phi^4$  local vertex by two (disconnected)  $\phi^2$  vertices (with derivative couplings). There is a loss of connectivity which prevents one from applying the key observation above in a too naive manner. It is here that phase space analysis becomes useful. The outline of the strategy is as follows. Either the graphs contributing to  $a_n$  or  $b_n$  are spread over a large number  $p$  of cubes of the series of scaled lattices naturally associated to the multiscale decomposition, or they are concentrated in a few such cubes. The transition is somewhat arbitrary but may be taken at  $p \simeq \frac{n}{(\log n)^\delta}$  for some small  $\delta > 0$ . In the first case, using the by now familiar horizontal and vertical decay associated to phase space (recall that such decay requires the use of “usefully

renormalized” expansions) we should prove that the corresponding contributions do not contribute at all to the leading behavior at large  $n$ . In the second case, one should think of the vertices as densely packed into small areas of phase space. In this case one can apply the Sobolev inequality; the loss of connectivity is harmless (as far as leading large order behavior is concerned) because the total volume is then small.

Returning to the language of Sect II.3, we introduce the triplets  $(G, \mathbf{F}, \mu)$  made of a graph, an assignment, and a forest  $\mathbf{F}$  which is safe for  $\mu$ , and the corresponding integrands

$$Z_{G,\mathbf{F},\mu} \equiv \prod_{g \in \mathbf{F}} (-\tau_g^*) \prod_{g \in \mathbf{H}_\mu(\mathbf{F})} (1 - \tau_g^*) Z_{G,\mu} \quad (\text{II.6.17})$$

where  $Z_{G,\mu}$  is as in (II.4.8) the ordinary Feynman integrand for the assignment  $\mu$ . Useful renormalization corresponds to the case  $\mathbf{F} = \emptyset$ , hence the bound (II.6.16) may be written as:

$$\sum_{(G,\mathbf{F},\mu); n(G)=n, \mathbf{F}=\emptyset} \int \prod_{v \in G} dx_v |Z_{(G,\mathbf{F},\mu)}| \leq n! [a_{Lip}]^n (1 + \epsilon(n))^n \quad (\text{II.6.18})$$

Next we decompose the (usefully renormalized) perturbation theory according to the multiscale slicing, as in Sect.II.3. The crucial point is to decompose the series in two according to whether the graphs and assignments  $(G, \mu)$  spread over a large region of phase space or not. To measure the size of a region in phase space, it is natural to consider that distance scales are the inverse of momentum scales. For each scale  $i$ ,  $i = 0, 1, 2, \dots$  this leads one to introduce the scaled lattice  $\mathbf{D}_i$  made of cubes of side  $M^{-i}$ , and define  $\mathbf{D} = \cup_i \mathbf{D}_i$  (see Fig.II.6.2). Now at a vertex  $v$  of a graph sitting at  $x_v$  there are at most 4 propagators which meet, with scales  $i_1(v), i_2(v), i_3(v)$  and  $i_4(v)$ , the maximum being  $e_v(\mu)$  by (II.4.4), and we may associate to  $v$  the set of the four cubes  $\Delta_1, \dots, \Delta_4$ ,  $\Delta_j \in \mathbf{D}_{i_j(v)}$  to which  $x_v$  belongs. If we repeat this for each vertex of a graph, we obtain a region which is the natural domain in  $\mathbf{D}$  of the contribution associated to  $(G, \mu)$ . This domain  $X_0^{vert}(G, \mu) \subset \mathbf{D}$  is simply obtained by coloring the cubes of  $\mathbf{D}$  which contain the meeting of a line at a vertex in our standard representation Fig.II.1.2 of phase space (see Fig.II.6.3). It will be called the “vertex domain” of the contribution in phase space, and a good measure of its size is simply the total number  $x_0^{vert}(G, \mu)$  of cubes in it.

However there is a problem with this simple approach. Remember that our strategy is to bound crudely the contributions spread over large domains of phase space (they should be small, anyway), and to apply the Sobolev inequality only to the sum over all graphs spread over a small domain of phase space. The Sobolev inequality (II.6.8) in a finite volume  $X_0$  applies with the same constant as in infinite volume only for fields in  $H_0^1(X_0)$ , the Sobolev space of functions  $\phi$  with square integrable gradients which vanish on the boundary of  $X_0$ , which is a natural subset of  $H_0^1(\mathbb{R}^4)$ . But the sum of graphs that we are studying is related to an integral  $(\int_{X_0} \phi^4)^n d\mu_C(\phi)$  where the gaussian measure  $d\mu_C(\phi)$  corresponds to the propagator  $C$ . Using this propagator creates a problem because the sample fields for  $d\mu_C$  are not in  $H_0^1(X_0)$ . This is not a regularity problem since in a finite region

of phase space the propagator has an ultraviolet cutoff of type (I.3.7) so that the sample fields are very smooth, but the problem is that these sample fields have no reason to vanish on  $\partial X_0$ . This would be however the case for instance if we could use a propagator in which the Laplacian entering the propagator's definition has Dirichlet boundary conditions on  $X_0$  (for the definition of gaussian measures and their support properties, see [Er1],[Si1]). For such a task it is convenient to introduce the Wiener path representation of the propagator [GJS],[FO]. Fortunately this representation is fully compatible with the parametric representation, hence with cutoffs (I.3.7) or (II.1.3), since  $\alpha$  is simply the proper time of the path. It is:

$$C^i(x, y) = \int_{M^{-2i}}^{M^{-2(i-1)}} dt e^{-m^2 t} \int P_t(x, y) d\omega \quad (\text{II.6.19})$$

where  $P_t(x, y) d\omega$  is the Wiener measure on the sets of all paths starting at  $x$  at time 0 and ending at  $y$  at time  $t$ .

Our goal of using the Sobolev inequality then leads us to consider a larger domain in phase space called the “vertex and propagator” domain. In [MNRS] this domain arises in a natural way as the result of an inductive cluster expansion of the Glimm-Jaffe Spencer type [GJS]. Since we postpone the definition of cluster expansions to the next section, we will give an equivalent global (set-theoretic) definition of this domain.

We decompose first each propagator in the perturbative expansion as in (II.6.19), so we rewrite the perturbative expansion as a sum over triplets  $(G, \mu, \Omega)$  where  $G$  and  $\mu$  are a graph and an assignment, as before, and  $\Omega$  is a set of  $l(G)$  paths  $\omega_1, \dots, \omega_{l(G)}$ ; the first two sums are discrete, but the sum over  $\Omega$  means a product of Wiener integrals  $\prod_{l=1}^{l(G)} \int P_{t_l}(x_l, y_l) d\omega_l$ . For this generalized multiscale representation, the straight lines which represent propagators in the standard picture (Fig.II.1.2) should be replaced by the more complicated paths of Fig.II.6.4.

Then it may seem natural to define the “full propagator domain” as the union over  $l$  of all the cubes of  $\mathbf{D}_{i_l(\mu)}$  visited by the path  $\omega_l$ . However this naive domain is again not the correct concept in this case, because we look for a definition such that when the propagator domain is large (compared to  $n/(\log n)^\delta$ ), the corresponding contributions are small (in the sense of Lemma II.6.2 below). This will not be the case with the above naive definition, for instance when the vertex domain is small and the “full propagator domain” is large only because of a *single* propagator whose path extends over a large set of cubes. In fact a definition adapted to our purpose is the following:

Let us consider a subset  $X = \{\Delta_1, \dots, \Delta_j, \dots, \Delta_q\}$  of  $\mathbf{D}$ . We say that this subset has the property HZ (for “horizontally connected”) with respect to  $(G, \mu, \Omega)$  if and only if for any  $\Delta_j \in X \cap \mathbf{D}_i$  (hence for any cube of  $X$  of scale  $i$ ) there is a *distinct* line  $l_j$  of the same scale ( $i_{l_j}(\mu) = i$ ) such that the corresponding path  $\omega_{l_j}$  visits  $\Delta_j$ . (The important fact is that all the lines  $l_j$  have to be different:  $l_j \neq l_{j'}$  for  $j \neq j'$ ).

Then we define the propagator domain  $X_0^{prop}(G, \mu, \Omega)$  of  $(G, \mu, \Omega)$  as one (arbitrarily chosen) subset of maximal cardinality  $x_0^{prop}(G, \mu, \Omega) = |X_0^{prop}(G, \mu, \Omega)|$  among all subsets which have the property HZ. Associated to this propagator domain, there is therefore a particular set of distinct lines  $l_j$ ,  $j = 1, \dots, x_0^{prop}(G, \mu, \Omega)$ , such that each  $\omega_{l_j}$  visits the corresponding cube  $\Delta_j$  of  $X_0^{prop}(G, \mu, \Omega)$ . We choose

a particular set of such lines,  $L(G, \mu, \Omega)$ , and call it the horizontal connections of the propagator domain. It is important to notice that for  $l \in L(G, \mu, \Omega)$  the path  $\omega_l$  does not necessarily stay in  $X_0^{prop}(G, \mu, \Omega)$ , but for  $l$  not in  $L(G, \mu, \Omega)$ , the path  $\omega_l$  must be confined in the cubes of  $X_0^{prop}(G, \mu, \Omega)$ ; otherwise  $x_0^{prop}(G, \mu, \Omega)$  would not be maximal (see Fig.II.6.4 for examples of the natural naive propagator domain and of the propagator domain as defined above).

That this definition is an appropriate one for our purpose will become clear below; but it is not the unique one possible\*, and it requires some arbitrary choices, which is unfortunate on an aesthetic level, but does not seem to be easily avoidable: arbitrary choices appear in most cluster expansions, often in the form of an arbitrary ordering of some finite set of geometric objects.

We define now:

$$X_0(G, \mu, \Omega) \equiv X_0^{vert}(G, \mu) \cup X_0^{prop(G, \mu, \Omega)} \quad \text{and} \quad x_0(G, \mu, \Omega) \equiv |X_0(G, \mu, \Omega)| \quad .$$

The proof of (II.6.18) is decomposed into two steps. Fix  $\delta > 0$ , and write:

$$a_n^{UR} = a_{n,small}^{UR} + a_{n,large}^{UR} \quad (\text{II.6.20})$$

where  $a_{n,small}$  is the sum over contributions with  $x_0(G, \mu, \Omega) \leq n(\log n)^{-\delta}$ , and  $a_{n,large}$  is the complement. We want to bound  $a_{n,large}^{UR}$ . We may again distinguish two subcases. Let  $\delta' > 0$  be such that  $\delta \ll \delta'$ . The first subcase is when  $x_0^{vert}(G, \mu) \leq n(\log n)^{-\delta'}$ , and we call the sum of the corresponding contributions  $a_{n,large \ prop}^{UR}$ . In this first subcase, we remark indeed that we must have  $x_0^{prop}(G, \mu, \Omega) > (n/2)(\log n)^{-\delta}$  (since  $\delta \ll \delta'$ ). In the second subcase, when  $x_0^{vert}(G, \mu) > n(\log n)^{-\delta'}$ , the sum of contributions is called  $a_{n,large \ vert}^{UR}$ .

Then our first goal is to explain, without entering all technicalities, the following bounds:

**Lemma II.6.1** For some  $\epsilon$  (depending on  $\delta$ ):

$$|a_{n,large \ prop}^{UR}| \leq C^n \frac{n!}{e^{n(\log n)^\epsilon}} \quad (\text{II.6.21a})$$

$$|a_{n,large \ vert}^{UR}| \leq C^n \frac{n!}{e^{n(\log n)^\epsilon}} \quad (\text{II.6.21b})$$

Let us start with (II.6.21a). The reason for which  $a_{n,large \ prop}^{UR}$  is small is that many propagators are longer than their typical decay scale. More precisely since there at most  $const.r^4$  cubes in  $\mathbf{D}_i$  at a distance less or equal to  $rM^{-i}$  of any given cube of  $X_0^{vert}(G, \mu)$ , there must be at least half of the cubes of the “propagator domain” which are far (in the relevant scale) of any cube of the “vertex domain” of the same scale. More precisely, at least half of the cubes  $\Delta$  of  $X_0^{prop}(G, \mu, \Omega)$  must satisfy an inequality:

$$\text{dist}(\Delta, X_0^{vert}(G, \mu) \cap \mathbf{D}_i) \geq cM^{-i}(\log n)^{\frac{\delta'-\delta}{4}} \quad \text{if } \Delta \in \mathbf{D}_i \quad (\text{II.6.22})$$

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\* To test its understanding of the problem we suggest that the reader tries to invent an other one, for instance with at most *two* cubes associated to distinct propagators.

The full propagator  $C^i(x, y)$  satisfies a bound like (II.1.6), but the conditioned propagator:

$$C_\Delta^i(x, y) = \int_{M^{-2i}}^{M^{-2(i-1)}} dt e^{-m^2 t} \int P_t(x, y) \chi_\Delta(\omega) d\omega \quad (\text{II.6.23})$$

where  $\chi_\Delta(\omega)$  is the characteristic function forcing  $\omega$  to visit  $\Delta$ , satisfies a more detailed bound:

$$C_\Delta^i(x, y) \leq K M^{2i} e^{-\eta M^i [\text{dist}(x, \Delta) + \text{dist}(\Delta, y)]} \quad (\text{II.6.24})$$

where  $\eta$  is a small number (the notation  $\delta$  like in (II.1.6) would be confusing).

This bound is intuitively obvious if we recall that the travel time on paths acts as an infrared cutoff. To prove it is a standard exercise in using the “additive” Markovian structure of the Wiener measure, and decomposing the path  $\omega$  into two pieces by introducing the first hitting time of  $\omega$  with  $\Delta$  (see e.g. [GJS]). (Actually, remark that in the literature it is standard to use as primary objects the faces which make the boundary of the cubes rather than the cubes themselves; this point of view helps to write a more systematic cluster expansion).

We may collect half of the decay (II.6.24) from the propagators of the horizontal connections  $L(G, \mu, \Omega)$  and combine it with the information (II.6.22) to extract a factor

$$\prod_{j=1}^{(n/4)(\log n)^{-\delta}} e^{-\eta(\log n)^{\frac{\delta'-\delta}{4}}} = e^{-\eta' n(\log n)^{\frac{\delta'-5\delta}{4}}} \quad (\text{II.6.25})$$

(Remark that here it is crucial to know that each cube of the propagator domain is associated to a distinct line of the graph).

We may take  $\delta' - 5\delta > 0$ . After extraction of the factor (II.6.25), propagators still satisfy the bound (II.1.6) (with a different constant for the decay). Hence (now forgetting the restrictions introduced by the condition *large prop*) we can apply Theorem II.3.1-2, and bound  $|a_{n, \text{large prop}}|$  by  $K^n n! e^{-n(\log n)^\epsilon}$  for some  $\epsilon > 0$ ; roughly speaking  $K^n$  comes from the estimate of Theorem II.3.1-2 for single amplitudes,  $n!$  comes from the number of graphs, and  $e^{-n(\log n)^\epsilon}$  from the factor (II.6.25). This establishes (II.6.21a).

The second subcase,  $a_{n, \text{large vert}}^{UR}$ , is bounded according to a slightly different idea: it is no longer the value of each individual amplitude which is small because of propagators longer than usual, but it is “statistically” that the sum of amplitudes is not as large as  $c^n n!$ , as should be expected from the number of Wick contractions divided by the vertex symmetry factor, whose leading behavior is  $\frac{4n!!}{n!} \simeq n!$ . Let us illustrate first on a simpler example this idea, which one may call the “local factorials” principle.

### Lemma II.6.2: the local factorial principle

Let us consider the  $\phi^4$  theory (in any dimension) with fixed cutoff  $\kappa = 1$  (so only the first slice  $i = 0$  is kept). Let  $n(\Delta)$  be a family of integers associated to each cube  $\Delta \in \mathbf{D}_0$  such that  $\sum_\Delta n(\Delta) = n$ . We call  $a_n^{\{n(\Delta)\}}$  the sum of the

contributions of the graphs which have exactly  $n(\Delta)$  vertices in each  $\Delta$ . There exists a constant  $K$  such that:

$$|a_n^{\{n(\Delta)\}}| \leq K^n \prod_{\Delta} n(\Delta)! \quad (\text{II.6.26})$$

The proof of this local factorial principle uses the exponential decay of the propagator, which suppresses strongly the Wick contractions which join distant cubes. The result has the same form (up to a large value of  $K$ ) as if in fact all Wick contractions between *distinct* cubes were suppressed, which means that the usual factorial  $n!$  reflecting the number of Wick contractions of  $4n$  fields divided by the vertex symmetry factor  $n!$  is replaced by a product of local such factorials in each cube. To prove the lemma we start with  $n$  vertices. The multinomial factor  $\frac{n!}{\prod_{\Delta} n(\Delta)!}$  allows to distribute them in the cubes. Then we build Wick contractions by choosing for a given field, first the cube containing the field to which it contracts, then the particular field to which it contracts, and iterating this process until all the fields are exhausted. The first choice (of the cubes) lead to a sequence of sums, each of which is controlled by the decay of the corresponding propagators, and leads to a constant per sum hence to  $c^n$ . The second choice (the field in the cube) leads to a factor  $c'^n \prod_{\Delta} n(\Delta)!^2$ , provided at each step  $j$ ,  $j=1,\dots,2n$ , we contract a field in a cube  $\Delta_j$  which contains a maximal number  $n_j(\Delta)$  of remaining fields not yet contracted; hence the choice of the field in  $\Delta'_j$  to which it will contract will cost a factor  $n'_j \leq \sqrt{n_j \cdot n'_j}$  and the total process will cost a factor bounded by  $\sqrt{\prod_{\Delta} n(\Delta)^{4n(\Delta)}}$  as announced. Taking the multinomial coefficient into account achieves the proof of the Lemma.

We remark that when the number  $p$  of occupied cubes gets large, the behavior of  $\prod n(\Delta)!$  may become significantly smaller than  $n!$ ; it is certainly bounded by  $\frac{n!}{p!}$ . This motivates the following stronger result, in which the domain of occupied cubes and the numbers  $n(\Delta)$  are no longer given, but summed up:

**Lemma II.6.3: Large “vertex domain” bound.**

Let  $p \leq n$  be an integer. For the one slice model, the sum of all perturbative contributions to a connected Schwinger function which contain vertices in exactly  $p$  different cubes,  $a_{n,p}$  satisfies:

$$a_{n,p} \leq K^n \frac{n!}{p!} \quad (\text{II.6.27})$$

To prove this result it is natural to use a (single-scale) cluster expansion like the one introduced in the next section (III.1). Roughly speaking, the convergence of the Brydges-Battle-Federbush tree cluster expansion, Theorem III.1.1, applied to this problem, means that we can choose the cubes containing the vertices and build a tree of Wick contractions connecting them at the cost of only  $c^p$  (hence without any factorial factor). This tree eats up  $2p$  fields, and the remaining Wick contractions create only a factor of order  $(4n - 2p)!! \simeq \frac{2n!}{p!}$ , which makes the result plausible. The true proof, however, is more complicated because there is some book-keeping of vertex symmetry factors to do, and we give it in section III.1, as an example of application of the (single slice) cluster expansion.

We understand now that when the contributions  $a_{n,large\ vert}^{UR}$  are restricted to the 0-th slice, the desired bound (II.6.21b) holds, since  $\frac{n!}{(n(\log n)-\delta')!} \leq c^n \frac{n!}{e^{n(\log n)\epsilon}}$  for some  $\epsilon$ . Hence to extend this bound to  $a_{n,large\ vert}^{UR}$  is the same thing than to extend the local factorial principle and lemma II.6.3 from a single slice model to the general phase space situation. This can be done because the vertical and horizontal exponential decay of usefully renormalized perturbation theory is the correct generalization to phase space of the horizontal “spatial decay” of the single slice model. However this requires the generalization to phase space of the single-slice (standard) cluster expansion (which is what we call a “multiscale cluster expansion”), and even with this tool the details of the combinatoric ([MNRS, section 3 and Appendix B) remain complicated. Since the next chapter of this book is devoted to cluster expansions and their use in constructive theory, we suggest that after studying this chapter, the interested reader returns to (II.6.21b) and builds up its own proof for it, using the simplest possible multiscale expansion (of the “pair of cubes” type) rather than the one of [GJS],[MNRS] (to use section 3 of [MNRS] for some clues is of course allowed!). Indeed for a proof of (II.6.21b) a cluster expansion which localizes only vertices is clearly enough; it is only for (II.6.21a) and Lemma II.6.5 below that a cluster expansion which localizes both vertices and propagators is required.

It remains now to apply the Sobolev inequality when the total vertex and propagator domain in phase space is small. We follow the same approach as for the large vertex domain, namely we start with an easy lemma as a motivation:

#### **Lemma II.6.4 Upper Lipatov bound in finite volume**

Let us consider the  $\phi_4^4$  model with fixed ultraviolet cutoff  $\kappa$ . Let  $X_0$  be a region of total volume  $|X_0| \leq \frac{n}{(\log n)^\delta}$  with  $0 < \delta < 1$ . There exists  $\epsilon(n)$  such that  $\lim_{n \rightarrow \infty} \epsilon(n) = 0$  and:

$$\frac{1}{n!} \int d\mu_\kappa(\phi) \left[ \int_{X_0} \phi^4(x) dx \right]^n \leq (1 + \epsilon(n))^n n! a_{Lip}^n \quad (\text{II.6.28})$$

where we recall that  $a_{Lip}$ , the Lipatov constant, is defined by (II.6.5), hence  $a = (4K_{Lip})^2$ ,  $K_{Lip}$  being defined by (II.6.7-8) for  $d = 4$ .

**Proof** Applying the Sobolev inequality (II.6.7) we obtain:

$$[\text{left hand side of (II.6.28)}] \leq \frac{1}{n!} K_{Lip}^{2n} \int d\mu_\kappa(\phi) \left[ \int_{X_0} \partial_\mu \phi \partial^\mu \phi + m^2 \phi^2 \right]^{2n} \quad (\text{II.6.29})$$

Integrating over  $d\mu_\kappa(\phi)$ , we get a sum over closed loops, each propagator being  $(-\Delta + m^2)C_\kappa$ , where  $C_\kappa$ , the cutoff propagator (I.3.7-8) is smaller than  $(p^2 + m^2)^{-1}$  in Fourier space and decays exponentially in direct space (over lengths of order  $\kappa^{-1}$ ). Therefore each closed loop contribution is bounded by  $C|X_0|$ , where  $C$  is a constant depending on  $\kappa$  (the single factor  $|X_0|$  takes into account the single translation invariance of the connected loop). Hence:

$$\int d\mu(\phi) \left[ \int_{X_0} \partial_\mu \phi \partial^\mu \phi + m^2 \phi^2 \right]^{2n}$$

$$\leq \sum_{p=1}^{2n} \frac{(2n)!}{p!} (C|X_0|)^p \sum_{\substack{t_1, \dots, t_p \\ \sum_{j=1}^p t_j = 2n}} \frac{(2t_1 - 1)!! \dots (2t_p - 1)!!}{t_1! \dots t_p!} \frac{(2n)!}{p!} (C|X_0|)^p 2^{2n} \sum_{p=1}^{2n} C_{2n}^p \quad (\text{II.6.30})$$

where  $C_{2n}^p$  is the binomial coefficient,  $p$  is the total number of closed loops, and  $t_1, \dots, t_p$  are the number of vertices in each loop. The combinatorial factors arise from the number of possible Wick contractions for each loop, which is easy to compute. Since  $C|X_0| < \frac{n}{(\log n)^\delta}$ , we have obviously:

$$\sum_{p=1}^{2n} C_{2n}^p \frac{(C|X_0|)^p}{p!} \leq e^{Cn(\log n)^{-\epsilon}} \quad \text{for some } \epsilon \quad (\text{II.6.31})$$

Remembering  $(2n)! \leq 2^{2n}(n!)^2$ , this proves the lemma.

It remains to extend this lemma to the phase space context:

### Lemma II.6.5

$$|a_{n,small}^{UR}| \leq n![a_{Lip}]^n (1 + \epsilon(n))^n \quad (\text{II.6.32})$$

Several difficulties arise when one tries to extend Lemma II.6.4 to Lemma II.6.5. First in Lemma II.6.4 the region  $X_0$  is given, but in Lemma II.6.5 we know only that it is small and we must sum over all possibilities. A similar difficulty separates Lemma II.6.2 from II.6.3; and we said that the solution requires a kind of cluster expansion. This is true also here, and this cluster expansion must be a multiscale one which localizes both propagators and vertices. Again here we provide simply minimal guiding remarks, in order to avoid overlap with the next chapter. A cluster expansion is, roughly speaking, an algebraic machinery to define a domain and select a particular explicit subset of connections between the different regions (here, cubes) of the domain, which ensure sufficient decay between these regions, so that one can sum upon the position and shape of the domain. Concretely this is usually done by Taylor expansions in interpolating (“decoupling”) parameters. But in perturbation theory (which is free of the stringent positivity requirements of constructive theory) these interpolating parameters are not strictly speaking necessary; in [MNRS] a mixed approach is used in which for the horizontal (in the sense of Fig.II.2.1) cluster expansions, interpolating parameters are used, but for the vertical connections, a set of vertices with fields in different slices (vertical dotted lines in our standard multiscale representation) is selected in a set-theoretic way (i.e. without interpolating parameters and Taylor expansions). Here since we define the full “vertex and propagator domain” in a set-theoretic way it is more natural to select the set of horizontal connections also in a set theoretic way, so that no interpolating parameters and Taylor formulae are ever needed; this set is simply  $L(G, \mu, \Omega)$ , defined as above. This simplifies some details in [MNRS].

All together this process selects particular propagators and vertices which connect together the vertex and propagator domain, and contain enough convergence factors to realize the vertical and horizontal exponential decay between its cubes, hence to solve the problem of summing over all such possible domains.

Roughly speaking, for convergent almost local subgraphs, it is enough to select vertices in them with at least 5 external fields of lower momenta. When divergent subgraphs of the almost local type appear in some places, the corresponding useful counterterms are associated to them. We may select all the external vertices of these subgraphs and let the corresponding subtractions act on their external lines. In both cases vertical decay is generated, the main problem being the burden of notations. The price to pay is simply that all fields hooked to selected vertices cannot be included in the set of remaining fields to which the Sobolev inequality is applied. However an important point is that, as a whole, the number of fields contracted into horizontal connections or hooked to selected vertices remain bounded by  $\text{const.}x_0(G, \mu, \Omega)$ , hence by  $\text{const.} \frac{n}{(\log n)^\delta}$ . Hence the remaining fields are vastly the majority.

There are still implicit restrictions on the Wick contractions of these remaining fields (for instance the ones who tell that  $G$  itself must be connected). We can now lift all these restrictions, (hence producing an overestimate for  $a_{n,small}^{UR}$ ). Then the remaining fields may be written as a single functional integral of the type  $\int (\int_{X_0} \phi^4)^{n'} d\mu_{X_0}(\phi)$  and *no longer* developped into Feynman graphs. Remark that we use the word “field” for simplicity, but since Feynman graphs rather than fields were our starting point, one should more precisely say that the remaining “pieces of Feynman graphs” are written as a single functional integral. Also our notations are very loose. In this integral  $\int (\int_{X_0} \phi^4)^{n'} d\mu_{X_0}(\phi)$ , for simplicity we did not decompose the field  $\phi$  as the sum of fields associated to each slice (using (II.1.5)) as should be done in fact. Also for simplicity a single notation  $X_0$  recalls the restrictions over the range of integration for vertices coming from the definition of the vertex domain, and the restrictions on the propagator, hence on the gaussian measure corresponding to the definition of the propagator domain. The important point to stress is that this definition is chosen so that although the paths of the propagators  $l_j$  of the horizontal connections do not necessarily lie in the “vertex and propagator domain”, all the paths of the other propagators do. This implies that the gaussian measure  $d\mu_{X_0}$  with which these remaining fields are integrated is supported on fields which have their support inside this domain and vanish, together with their gradient, at the boundary  $\partial X_0$ . Hence these fields do obey the (infinite volume) Sobolev inequality. Therefore we can apply the Sobolev inequality in the manner of Lemma II.6.4 to this functional integral. Since the number of cubes in phase space turns out to be, as expected intuitively, the correct factor which generalizes the volume  $|X_0|$  in Lemma II.6.4, Lemma II.6.5 follows. This completes our sketch of the proof of Theorem II.6.3.

Finally we will sketch how the Lipatov upper bound on usefully renormalized series combines with an analysis of the recursion relation for effective coupling constant in order to obtain Theorem II.6.2 [DFR]. We will need in fact the slightly more detailed corollaries of the analysis above:

For any  $\delta > 0$  there exists a function  $\epsilon(n)$  with  $\lim_{n \rightarrow \infty} = 0$  and

$$\sum_{(G, \mathbf{F}, \mu); n(G)=n, \mathbf{F}=\emptyset} \int \prod_{v \in G} dx_v |Z_{(G, \mathbf{F}, \mu)}| e^{(2-\delta)i_{max}(\mu)} \leq n! [a_{Lip}]^n (1 + \epsilon(n))^n \quad (\text{II.6.33})$$

where  $i_{max}(\mu)$  is the maximum index in  $\mu$ . This bound will result from the fact that

in index space the exponential decay rate is at least 2, since the worst superficial convergence degree is at least 2 after useful renormalization has been performed. After the factor  $e^{(2-\delta)i_{max}(\mu)}$  of (II.6.33) is included, there remains indeed a vertical exponential decay with rate at least  $\delta$ .

We need also related bounds for the 4 point counterterms:

$$\sum_{\substack{(G,\mathbf{F},\mu) \\ n(G)=n, N(G)=4, \mathbf{F}=\emptyset \\ k \leq i_G(\mathbf{F}) \leq e}} \int \prod_{v \in G} dx_v |Z_{(G,\mathbf{F},\mu)}^0| e^{(2-\delta)[i_{max}(\mu)-i_G(\mathbf{F})]} \\ \leq (e-k+1)n![a_{Lip}]^n(1+\epsilon(n))^n \quad (\text{II.6.34})$$

where  $Z^0$  is the coupling constant counterterm obtained by applying  $\tau_G$  instead of  $(1-\tau_G)$  in (II.6.17). Remark that now we have only exponential decay between the maximal index  $i_{max}(\mu)$  and the minimal index  $i_G(\mathbf{F})$  (as defined by (II.3.34)) in  $G$ . There is one global translation invariance in vertical index space for  $G$  corresponding to the logarithmic divergence of a coupling constant counterterm. These features explain the form of (II.6.34).

Similarly:

$$\sum_{\substack{(G,\mathbf{F},\mu) \\ n(G)=n, N(G)=2, \mathbf{F}=\emptyset \\ k \leq i_G(\mathbf{F}) \leq e}} \int \prod_{v \in G} dx_v |Z_{(G,\mathbf{F},\mu)}^0| e^{(2-\delta)[i_{max}(\mu)-i_G(\mathbf{F})]} \\ \leq M^{2(e-k)}n![a_{Lip}]^n(1+\epsilon(n))^n \quad (\text{II.6.35})$$

where  $Z^0$  is the mass counterterm obtained by applying  $\tau_G^0$  instead of  $(1-\tau_G)$  in (II.6.17). This bound reflects the quadratic divergence of a mass counterterm (which in practice in the case of useless counterterms is always compensated by the quadratic convergence of one of the external legs of  $G$ ).

In [DFR] the renormalized perturbation series is then recast into a form which is not exactly the effective expansion of Sect.II.4, but an expansion intermediate between this one and the renormalized expansion. More precisely only the useless counterterms for the subgraphs isomorphic to the bubble are resummed into effective constants, and the other useless counterterms for all subgraphs except the ones isomorphic to the bubble are kept in the expansion. This may seem complicated and not very natural, but it has the advantage of leading to a very simple recursion rule for effective constants which is simply a second order polynomial [DFR]:

$$g^i(g_r) = g^{i-1}(g_r) + Y(i-1)(g^{i-1}(g_r))^2 \quad (\text{II.6.36})$$

where  $Y(i)$  is the value of the bubble graph with minimal index  $i$ :

$$Y(i) = b \sum_{\min\{j_1, j_2\}=i} \int d^4x C^{j_1}(0, x) C^{j_2}(x, 0) \quad (\text{II.6.37})$$

$b$  being the symmetry factor of the bubble. The advantage is that this simple recursion leads to an easy bound on the Borel transform of the product of effective couplings associated to a contribution  $(G, \mu, \mathbf{F})$ . This product  $DR_{(G, \mu, \mathbf{F})} =$

$\prod_{v \in G} g^{e_v(\mu, \mathbf{F})}$ , using (II.6.36-37) satisfies indeed:

$$\begin{aligned} |B'(DR_{(G,\mu,\mathbf{F})})(b)| &\leq \frac{|b|}{(n-1)!} e^{\sum_{0 \leq j < i_{max}(\mu)} Y(j)|b|} \\ &\leq \frac{|b|^{n-1}}{(n-1)!} e^{|b|\beta_2 i_{max}(\mu) + const.} \end{aligned} \quad (\text{II.6.38})$$

where  $\beta_2 = 9/2\pi^2$  is the coefficient in Theorem II.6.2, and the modified Borel transform (easily related to the ordinary one) is simply  $B'(\sum_{n \geq 1} a_n g_r^n) \equiv \sum_{n \geq 1} \frac{a_n b^n}{(n-1)!}$ . Then in [DFR] some analysis of the useless counterterms is performed by combining the bounds (II.6.34-35) with some tedious combinatoric arguments. The outcome is that as long as  $|b| \leq 2/\beta_2$  the bound (II.6.38) essentially tells that there is still vertical index space decay and (II.6.33) leads to analyticity in the Borel plane, because the Lipatov constant is bigger than  $2/\beta_2$  by a factor 3/2. Theorem II.6.2 then follows.

The use of a modified effective expansion is a technical device not very appealing. Also Theorem II.6.2 remains unsatisfactory. We would like to rule out possible (“miraculous”) cancellations and have a proof of the existence of the first renormalon singularity at  $b = \beta_2/2$ , by showing that some derivative of  $B(b)$ , the Borel transform of the renormalized series blows up at this point. This result would be the first direct proof that the  $\phi_4^4$  series diverge and are not Borel summable. Ecalle’s theory of resurgent functions [Ec] (see also [GKT]) may be relevant for this difficult problem, because it gives information on the asymptotic behavior of the Borel transform of recursive relations much more general than (II.6.36). However the problem is difficult, because the natural approach would be to define rigorously the full  $\beta$  function. It is not yet clear whether this is possible in the BPHZ scheme [Kop]\* From [FMRS5] we know that the  $\beta$  function for infrared  $\phi_4^4$  is Borel summable in terms of the *bare* coupling; transcribed to the ultraviolet problem these results only mean that the  $\beta$  function of the massive theory with an ultraviolet cutoff  $M^\rho$  is Borel summable in terms of the renormalized coupling in a disk which shrinks as  $\rho^{-1}$  as  $\rho \rightarrow \infty$ . If we cannot use the full ultraviolet  $\beta$  function the only available approach seems to improve on the explicit resummations of [DFR]; if they could be extended to third order graphs, so as to reconstruct the first two terms of the  $\beta$  function rather than the first one, a proof of existence of the first  $\phi_4^4$  renormalon might become possible, in particular in the case of a vector model with large number of components, in which remainder terms might be bounded in  $\frac{1}{N}$  so as to exclude the possibility of miraculous cancellations.

This completes our review of the mathematical problem of  $\phi_4^4$  large order behavior. Although our guided tour of [MNRS]-[DFR] is no mathematical substitute for more rigorous proofs, we hope that it may provide some help in understanding Theorem II.6.2. Also we think that a study of this problem, in particular of

\* In the minimal subtraction scheme based on dimensional regularization there is some belief that the  $\beta$  function should be Borel summable [BDZ], but a proof may have to wait until constructive theory finds a non-perturbative way to define dimensional regularization...

the Lipatov upper bound, provides a natural introduction to the main theme of constructive theory, hence to the next part of this book.

Indeed although at first sight the large order behavior of  $\phi^4$  is a purely perturbative problem, we see that for this problem the language of Feynman graphs reaches its limits, and functional integration becomes unavoidable, precisely at the point where the Sobolev inequality is used to gain information on the collective behavior of the graphs. Very similar phenomena occur in constructive theory. The main theme in constructive theory, as we shall see now, is indeed exactly to sort out explicitly some critical connections between regions of phase space (by a cluster expansion), but to treat the rest of the theory as a functional integral. In constructive theory the goal is no longer to apply a Sobolev inequality to this rest, but it is still to take the functional form into account to gain some information which is hidden at the level of Feynman graphs, for instance the positivity used in “domination” or large fields bounds.

**PART III: CONSTRUCTIVE RENORMALIZATION**

### III.1 Single scale cluster and Mayer expansions

Cluster and Mayer expansions are key tools in many areas of mathematical physics. Introduced in constructive field theory by Glimm, Jaffe and Spencer to complete the construction of  $\phi_2^4$  [GJS], they have been improved or generalized over the years, in particular by Brydges, Battle and Federbush [BrFe][BaF1-3][Bat]. Unfortunately these expansions had for a while a reputation of being heavy to handle. In this section we try to dispel this impression by underlying the main facts behind their convergence. We do not try to present the best techniques or the optimal bounds. Our goal is simply to introduce beginners to the paradise of expansions; for a more complete review we suggest the reading of [Bry].

In this section we work always in a single momentum slice, for instance the slice with index 0 in phase space. In other words we have both fixed infrared and ultraviolet cutoff. Nevertheless we have in mind to use the single scale expansion as the building block for the multiscale expansions of the next section. This gives us some flexibility which we would not have if we were forced to treat as far as possible models with a single scale expansion. For instance, up to some extent the way the slicing is done, hence the form of these cutoffs may be chosen at will. In particular we can always require the sliced propagator to have fast decay at infinity, for instance exponential decay.

Taking advantage of this simplifies some arguments. For instance once the propagator decays sufficiently fast, we can use what we call the “volume argument”. This means simply that in a *finite dimensional* lattice of cubes, there are not many cubes close to a given one. This simple observation, combined with the rapid decay of the propagator, simplifies often the combinatoric which arises within the expansion.

We describe first how to recast a gaussian measure perturbed by a small stable interaction in the form of a polymer system with hard core interaction: the very name of polymer comes from this “excluded volume” effect, hence from the hard core interaction. This step is what we specifically call the cluster expansion. Then we show how to remove the hard core interaction and compute normalized quantities (or the pressure of the system). This is what we call the Mayer expansion, and it allows to control the thermodynamic (or infinite volume) limit. This Mayer expansion is really a systematic way to compute the logarithm of a grand canonical partition function. It would be an interesting exercise, which we leave to the reader, to apply this expansion formalism to ordinary perturbation theory (which may be considered as a grand canonical Bose gas of vertices). With this formalism one could for instance recover the reasons for which connected functions are sums of connected graphs and for which the combinatoric of the Bogoliubov recursion works.

The typical situation we study is a massive gaussian measure  $d\mu(\phi)$  with an ultra violet cutoff, perturbed by a  $g\phi^4$  interaction with  $g$  small. The measure could also be massless but with both infra red and ultra violet cutoffs (remember that it should correspond to a momentum *slice*.) This seems a somewhat trivial situation, but remember that even in the finite volume the perturbation expansion would diverge because of the large number of graphs. Moreover, even for a massive theory defined in a finite volume  $\Lambda$ , the thermodynamic limit  $\Lambda \rightarrow \infty$  is not trivial to define directly, and that is precisely what a single cluster and Mayer expansion

will easily do (the problem of the large number of graphs is bypassed because the expansion automatically develops only a piece (typically small) of the interaction).

The partition function in a volume  $\Lambda$  is

$$Z(\Lambda) = \int d\mu(\phi) e^{-g \int_{\Lambda} \phi^4(x) dx} \quad (\text{III.1.1})$$

the pressure is

$$P(\Lambda) = \frac{1}{\Lambda} \log Z(\Lambda) \quad (\text{III.1.2})$$

the unnormalized Schwinger functions are

$$S_{N,\Lambda}^u(z_1, \dots, z_N) = \int \phi(z_1) \dots \phi(z_N) d\mu(\phi) e^{-g \int_{\Lambda} \phi^4(x) dx} \quad (\text{III.1.3})$$

(remember that they are distributions which should be smeared by test functions  $\chi(z_1), \dots, \chi(z_N)$ ). The normalized Schwinger functions are

$$S_{N,\Lambda}(z_1, \dots, z_N) = \frac{1}{Z(\Lambda)} S_{N,\Lambda}^u(z_1, \dots, z_N) \quad (\text{III.1.4})$$

The interaction  $e^{-g \int_{\Lambda} \phi^4}$  could be generalized, but we should then assume that it remains stable, that is has a small constant like  $g$  in front, and that it is local in the sense that if we cut  $\Lambda$  in several regions  $\Lambda_i$ , it factorizes as a product of functions of the field in  $\Lambda_i$ : in our case this is just  $e^{-g \int_{\Lambda} \phi^4} = \prod_i e^{-g \int_{\Lambda_i} \phi^4}$ .

The dimension  $d$  of space time has to be finite (for the volume argument below) but is not necessarily 4 at this stage.

### A) The cluster expansion

We consider a normalized gaussian measure  $d\mu$  whose covariance is a symmetric positive definite operator in  $x$ -space  $C(x, y)$  which has good decrease at infinity. It could be for instance the covariance  $C^0$  (first slice) of (II.1.4), but in this section we drop the superscript 0 for simplicity. We assume that its decay is either exponential:

$$|C(x, y)| \leq O(1) e^{-|x-y|} \quad (\text{III.1.5})$$

or power-law with a very large (adjustable) power rate  $r$ :

$$|C(x, y)| \leq O(1) \left( \frac{1}{1 + |x - y|} \right)^r \quad (\text{III.1.6})$$

(We have rescaled the mass or the unit of length so that there is no scale coefficient in (III.1.5) or (III.1.6)). Then we divide our volume  $\Lambda$  into (hyper) cubes of side size unity, which form a lattice  $\mathbf{D}$ . On  $\mathbf{D}$  we adopt for convenience an arbitrary order:  $\mathbf{D} = \{\Delta^1, \dots, \Delta^{|\Lambda|}\}$ . We adopt also the same notation  $\Delta$  for a cube and for its characteristic function (also called sharp characteristic functions):  $\Delta(x) = 0$  if  $x \notin \Delta$ ,  $\Delta(x) = 1$  if  $x \in \Delta$ . We can write

$$\Lambda(x) = \sum_{\Delta \in \mathbf{D}} \Delta(x) \quad (\text{III.1.7})$$

and we might also use  $C_0^\infty$  versions of the characteristic functions  $\Delta$  (which we call smooth characteristic functions) so that (III.1.7) would then be a smooth partition of unity over  $\Lambda$ . In both cases the notation  $\int_\Delta f(x)$  means really  $\int_\Lambda \Delta(x)f(x)$ . Working within volume  $\Lambda$  means that we consider  $\Lambda(x)C(x,y)\Lambda(y)$  as our covariance.

Since coupling between different cubes, which prevents factorization of (III.1.1) over the cubes of  $\mathbf{D}$ , comes solely from  $d\mu$ , hence from the covariance  $C$ , we want to interpolate directly in  $C$  between an uncoupled situation and the coupled one. But we would like the interpolated covariance to preserve the positivity of  $C$  (as an operator) so that it still corresponds to an interpolated gaussian measure. To weaken off-diagonal elements (which generate the unwanted couplings) in a discrete, finite dimensional positive symmetric matrix  $C$ , one could multiply these off-diagonal elements  $C_{ij}$  and  $C_{ji}$  by an interpolating parameter  $s_{ij}$  for each pair  $i, j$ ,  $i \neq j$ . This point of view leads naturally to what is probably the simplest cluster expansion, the “pair of cubes” cluster expansion of Fig.III.1.1. However this process does not preserve positivity in general.

For instance  $\begin{pmatrix} 1 & 1 & 1 \\ 1 & 2 & 2 \\ 1 & 2 & 3 \end{pmatrix}$  is positive, but  $\begin{pmatrix} 1 & 1 & 0 \\ 1 & 2 & 2 \\ 0 & 2 & 3 \end{pmatrix}$  is not! Multiplying off-diagonal elements by parameters between 0 and 1 preserves positivity only if the diagonal elements are quite big compared to the off-diagonal ones. Hence to be used the pair of cubes expansion may require that we take a lattice of rather *large* cubes compared to the unit scale. This in effect reinforces the diagonal piece of  $C$  viewed as a matrix between the cubes of  $\mathbf{D}$  and allows preservation of positivity (see [FMRS5]). For the moment we do not want to use the (somewhat complicated) trick to enlarge our cubes. Then we can remark that in the particular case of a 2 by 2 matrix, positivity is preserved by damping the (single!) off-diagonal piece.

This suggests an inductive expansion in which at each step one tests the coupling between a set of cubes and the complement. This automatically preserves positivity and generates the tree-like expansion of Fig.III.1.2. It is this powerful point of view that we explain now; its only drawback is that it is inductive, hence its outcome is difficult to capture in a single formula.

For the tree expansion, we start with the first labelled cube in  $\mathbf{D}$ , say  $\Delta_1 = \Delta^1$  (the upper index is reserved for the arbitrary order on  $\mathbf{D}$ ), and we introduce a first parameter  $s_1$  which tests the coupling between  $\Delta_1$  and the rest of  $\mathbf{D}$  noted  $\bar{\Delta}_1$  (a bar indicates the complement in  $\mathbf{D}$ )

$$\begin{aligned} C(s_1, x, y) &= s_1 C(x, y) + (1 - s_1)[\Delta_1(x)C(x, y)\Delta_1(y) + \bar{\Delta}_1(x)C(x, y)\bar{\Delta}_1(y)] \\ &= \Delta_1 C \Delta_1 + \bar{\Delta}_1 C \bar{\Delta}_1 + s_1 [\Delta_1 C \bar{\Delta}_1 + \bar{\Delta}_1 C \Delta_1] \end{aligned} \quad (\text{III.1.8})$$

The first expression is a convex combination of functions of positive type, hence is of positive type, and proves that there is no problem to define the corresponding gaussian measure  $d\mu_{s_1}$ ; the second form is also useful to check that the  $s$  dependence is indeed on the off-diagonal terms. We follow this interpolation by a first order Taylor expansion. For instance for the partition function we write:

$$Z(\Lambda) = \int d\mu_{s_1}(\phi) e^{-g \int_\Lambda \phi^4} \Big|_{s_1=1}$$

$$\begin{aligned}
&= \int d\mu_{s_1} e^{g \int_\Lambda \phi^4} \Big|_{s_1=0} + \int_0^1 \frac{d}{ds_1} d\mu_{s_1}(\phi) e^{-g \int_\Lambda \phi^4} \\
&= Z(\Delta_1)Z(\bar{\Delta}_1) + \sum_{\Delta_2 \neq \Delta_1} Z_{\Delta_1 \Delta_2}(\Lambda)
\end{aligned} \tag{III.1.9}$$

with

$$Z_{\Delta_1 \Delta_2}(\Lambda) = \int_0^1 ds_1 \int d\mu_{s_1}(\phi) \int_{\Delta_1} dx \int_{\Delta_2} dy C(x, y) \left[ \frac{\delta}{\delta \phi(x)} \frac{\delta}{\delta \phi(y)} e^{-g \int_\Lambda \phi^4} \right] \tag{III.1.10}$$

(here it is necessary to distinguish between upper and lower indices because in general  $\Delta_2$  is not  $\Delta^2$ , the second cube in our list!). To check formulae (III.1.9)-(III.1.10) we recall that at  $s_1 = 0$  we can consider  $\Delta_1 \phi$  and  $\bar{\Delta}_1 \phi$  as two independent fields, and since the interaction factorizes,  $Z$  also factorizes. The advantage here is that we do not need really to worry about boundary conditions like Dirichlet, Neumann, etc.... which were important technicalities in former formalisms. Formula (III.1.10) follows from

$$\frac{d}{ds_1} C(s, x, y) = \sum_{\Delta_2 \neq \Delta_1} [\Delta_1(x)C(x, y)\Delta_2(y) + \Delta_2(x)C(x, y)\Delta_1(y)] \tag{III.1.11}$$

and from integration by parts with respect to  $\phi$ \*. It creates both a propagator between points localized in  $\Delta_1$  and  $\Delta_2$ , whose decay becomes available to perform the sum over  $\Delta_2$  at  $\Delta_1$  fixed in (III.1.9), and vertices hooked to this propagator, which are created by the derivatives  $\frac{\delta}{\delta \phi(x)}$  and  $\frac{\delta}{\delta \phi(y)}$  acting in (III.1.10) on the exponential. This is also useful because these vertices have precious factors  $g$  attached to them.

The definition of our second interpolation depends on which term in (III.1.9) we look at, and that's why this expansion is inductive. For the first term,  $\Delta_1$  is decoupled from the rest and we should start again with  $\mathbf{D}$  now restricted to  $\bar{\Delta}_1$ , hence pick the next cube (which, this time, is  $\Delta^2$ ) and test its coupling with  $\overline{\Delta_1 \cup \Delta^2}$ . But for the  $Z_{\Delta_1 \Delta_2}$  term we will introduce an interpolating parameter  $s_2$  which tests the coupling of the *union* of  $\Delta_1$  and  $\Delta_2$  to  $\overline{\Delta_1 \cup \Delta_2}$  (remember the former remarks about positivity). Hence in this case the interpolated propagator is a natural iteration of (III.1.8). Writing  $\Delta_{12}$  for  $\Delta_1 \cup \Delta_2$ ,  $\bar{\Delta}_{12}$  for its complement in  $\mathbf{D}$ , and using  $\Delta_1 \Delta_{12} = \Delta_1$ ,  $\Delta_1 \bar{\Delta}_{12} = 0$ , etc..., we get:

$$\begin{aligned}
C(s_1, s_2, x, y) &= s_1 s_2 C + s_1(1-s_2)[\Delta_{12} C \Delta_{12} + \bar{\Delta}_{12} C \bar{\Delta}_{12}] \\
&+ (1-s_1)s_2[\Delta_1 C \Delta_1 + \bar{\Delta}_1 C \bar{\Delta}_1] + (1-s_1)(1-s_2)[\Delta_1 C \Delta_1 + \Delta_2 C \Delta_2 + \bar{\Delta}_{12} C \bar{\Delta}_{12}] \\
&= \Delta_1 C \Delta_1 + \Delta_2 C \Delta_2 + \bar{\Delta}_{12} C \bar{\Delta}_{12} + s_1[\Delta_1 C \Delta_2 + \Delta_2 C \Delta_1] \\
&+ s_2[\Delta_2 C \bar{\Delta}_{12} + \bar{\Delta}_{12} C \Delta_2] + s_1 s_2[\Delta_1 C \bar{\Delta}_{12} + \bar{\Delta}_{12} C \Delta_1]
\end{aligned} \tag{III.1.12}$$

and the Taylor expansion in  $s_2$  gives:

$$Z_{\Delta_1 \Delta_2}(\Lambda) = Z_{\Delta_1 \Delta_2}(\Delta_{12})Z(\bar{\Delta}_{12})$$

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\* This functional integration by parts can be checked first for a polynomial integrand using the explicit rules of gaussian integration for this case; it can be extended to more general  $C^\infty$  functions by a continuity argument.

$$\begin{aligned}
& + \sum_{\Delta_3 \neq \Delta_1, \Delta_2} \int_0^1 ds_1 \int_0^1 ds_2 \int d\mu_{s_1, s_2}(\phi) \int_{\Delta_1} dx_1 \int_{\Delta_2} dy_1 \int_{\Delta_{12}} dx_2 \int_{\Delta_3} dy_2 \\
& \cdots C(x_1, y_1) C(s_1, x_2, y_2) \left[ \frac{\delta}{\delta\phi(x_2)} \frac{\delta}{\delta\phi(y_2)} \frac{\delta}{\delta\phi(x_1)} \frac{\delta}{\delta\phi(y_1)} e^{-g \int_{\Lambda} \phi^4} \right] \quad (\text{III.1.13})
\end{aligned}$$

If we want to have spatial integrals over unit cubes, we see that we need to break (III.1.13) according to whether  $x_2$  is in  $\Delta_1$  or in  $\Delta_2$ , so we get a term indexed by the set  $\Delta_1, \Delta_2, \Delta_3$ , and by the two possible trees  $\{(\Delta_1 \Delta_2), (\Delta_1 \Delta_3)\}$  or  $\{(\Delta_1 \Delta_2), (\Delta_2 \Delta_3)\}$  corresponding to the first and the second case respectively. We can iterate this process until the finite set  $\mathbf{D}$  is exhausted. Clearly the factorized contributions that we obtain, once they are decomposed so that each spatial integral is over a single cube, are not only indexed by their content or support in terms of cubes, but also by all possible tree structures that can be built on them and generated by the expansion. After reading part I, we should not be afraid anymore by trees anyway!

For summation we use a formalism of ordered sequences rather than sets because the former allows more easily the repetition of a single element. We should remark that in formula (III.1.9) for the first term  $Z(\Delta_1)Z(\bar{\Delta}_1)$  there are no  $\frac{\delta}{\delta\phi}$ 's, so no small constant generated. In fact the partition function of a single cube  $Z(\Delta_1) = Z(\Delta)$  is simply  $\int d\mu_{\Delta}(\phi) e^{-g \int_{\Lambda} \phi^4}$ , where  $d\mu_{\Delta}$  is the normalized gaussian measure corresponding to the covariance  $\Delta(x)C(x, y)\Delta(y)$ , and we expect  $Z(\Delta)$  to be close to 1, since  $g$  is small. To have really 1 we could do a first order Taylor expansion in  $g$  but it would complicate the picture. We prefer to compute a slightly different partition function:

$$\tilde{Z}(\Lambda) = Z(\Lambda)/Z(\Delta)^{|\Lambda|} \quad (\text{III.1.14})$$

( $|\Lambda|$  is of course the number of cubes in  $\Lambda$ , also equal to its volume since the cubes are of unit size). This corresponds to a finite shift of the pressure  $p(\Lambda)$  by the constant  $\log Z(\Delta)$ . With this choice, isolated cubes automatically cancel against the corresponding denominators in (III.1.14). The result of the cluster expansion is then simply:

$$\tilde{Z}(\Lambda) = \sum_{\substack{Y_1, \dots, Y_q \subset \mathbf{D} \\ Y_1, \dots, Y_q \text{ disjoint}}} \frac{1}{q!} A(Y_1) \cdots A(Y_N) \quad (\text{III.1.15})$$

where  $Y_1, \dots, Y_q$  are polymers (called also vacuum- or  $\emptyset$ -polymers), i.e., sets of at least 2 cubes of  $\mathbf{D}$  (isolated cubes or “monomers” are excluded), and the amplitude corresponding to a polymer  $Y = \{\Delta_1, \Delta_2, \dots, \Delta_p\}$  (where  $\Delta_1$  is the first cube of  $Y$  in the arbitrary order in  $\mathbf{D}$  and  $p = |Y|$ ) is:

$$\begin{aligned}
A(Y) = & \sum_{T'} \frac{1}{Z(\Delta)^{|Y|}} \int_0^1 \int_0^1 ds_1 \dots ds_{p-1} M_{T'}(s) \int d\mu_{\{s\}}(\phi) \left\{ \left[ \prod_{k=1}^{p-1} \int_{\Delta_{i_k}} dx_k \int_{\Delta_{j_k}} dy_k \right. \right. \\
& \left. \left. C(x_k, y_k) \frac{\delta}{\delta\phi(x_k)} \frac{\delta}{\delta\phi(y_k)} \right] e^{-g \int_Y \phi^4(x) dx} \right\} \quad (\text{III.1.16})
\end{aligned}$$

where to be precise, the sum is over ordered rooted trees  $T'$  with root 1, in the sense of Section I.4. Recall that these trees are ordered sequences of links  $l_k \equiv$

$(\Delta_{i_k}, \Delta_{j_{k+1}})$   $k = 1, \dots, p-1$ , connecting all the cubes of  $Y$  such that if we define  $\Delta_{j_1} \equiv \Delta_1$  we have, for any  $k \geq 1$ ,  $\Delta_{i_k} \in \{\Delta_{j_1}, \Delta_{j_2}, \dots, \Delta_{j_k}\}$  (see Fig.I.4.6). This ensures that for any  $k$  with  $1 \leq k \leq p$ ,  $Y_k(T') = \{\Delta_{j_r}, 1 \leq r \leq k\}$  has exactly  $k$  distinct cubes which are connected together by the subset  $T'_{k-1} = \{l_1, \dots, l_{k-1}\}$  of the  $k-1$  first links of  $T'$ ; in particular  $Y_p(T') = Y$ . As in Sect.I.4 it is convenient to consider the ancestor function  $a$  associated to  $T'$  such that the index  $a(k) \leq k$  is the one such that  $\Delta_{i_k} = \Delta_{j_{a(k)}}$ . The gaussian measure  $d\mu_{\{s\}}$  corresponds to the symmetric propagator  $C(s_1, \dots, s_{p-1}, x, y)$  which is simply  $C(x, y) \prod_{j=k'}^{k-1} s_j$  when  $y$  belongs to  $\Delta_{j_k}$  and  $x$  belongs to  $\Delta_{j_{k'}}$ , with  $k' \leq k$ ; it is zero when  $y$  or  $x$  are out of  $Y$ .

At an intermediate stage of the expansion, we may define

$$C(s_1, \dots, s_k, x, y) = C(s_1, \dots, s_{p-1}, x, y) \Big|_{s_{k+1} = \dots = s_{p-1} = 1}$$

There is therefore a polynomial  $s$  dependence of the propagators derived by the inductive expansion; by definition this polynomial dependence has been gathered in the factor  $M_{T'}(s)$ . It is easy to check that this factor is nothing but:

$$M_{T'}(s) = \prod_{k=1}^{p-1} \prod_{j=a(k)}^{k-1} s_j \quad (\text{III.1.17})$$

(III.1.15) expresses the partition function as the one of a gas of polymers swirling in  $\Lambda$ , with hardcore interaction (the condition of disjointness). From the decay of  $C$  and the smallness of  $g$  we expect this gas to be dilute, and typical polymers to be small.

When external fields  $\phi(z_1) \cdots \phi(z_N)$  are present we introduce the notion of an external polymer, which is a set  $Y$  of  $p$  cubes of  $\mathbf{D}$  together with a subset  $\zeta \subset \{z_1 \cdots z_N\}$  of external variables; it is then more precisely called a  $\zeta$ -polymer. Typically these external variables  $z_j$  are smeared against test functions  $\chi_j(z_j)$ , and one should keep in mind that if the support of these smearing functions does not overlap with the cubes of  $Y$  the corresponding amplitude of the polymer will be zero. Remark that we do not require  $\zeta$  to contain all the external variables which are localized in  $Y$  (i.e. have smearing functions with support in the cubes of  $Y$ ), a fact which will be convenient later.

Again we forbid the trivial case of a single isolated *empty* cube  $p = 1$ ,  $\zeta = \emptyset$ , but  $p = 1$  and  $\zeta \neq \emptyset$  is allowed. Then the truncated unnormalized Schwinger functions  $\widetilde{S}_{N,\Lambda}^u(z_1, \dots, z_N) = Z(\Delta)^{-|\Lambda|} S_{N,\Lambda}^u(z_1, \dots, z_N)$  are also expressed by the cluster expansion as a gas of polymers with hardcore interaction:

$$\widetilde{S}_{N,\Lambda}^u(z_1, \dots, z_N) = \sum_{\substack{Y_1, \dots, Y_q \subset \mathbf{D} \\ Y_1, \dots, Y_q \text{ disjoint linear in } \{z\}}} \frac{1}{q!} A(Y_1) \cdots A(Y_q) \quad (\text{III.1.18})$$

where the condition of “linearity in the external variables” means by definition that for each  $z_i$ ,  $1 \leq i \leq N$  there is exactly one polymer  $Y_j$ ,  $j \in [1, \dots, q]$  which contains it. The amplitude  $A(Y)$  is exactly similar to (III.1.16) but with  $e^{-g \int_Y \phi(x) dx}$

replaced by  $\left[ \prod_{z_j \in \zeta} \int \chi_j(z_j) dz_j \phi(z_j) e^{-g \int_Y \phi^4(x) dx} \right]$  if  $Y$  is a  $\zeta$ -polymer. The next step in the analysis is to perform a Mayer expansion that will take care of the hardcore interaction in (III.1.15) or (III.1.18). Anticipating on what follows, let us state that this expansion has been shown to converge, for such a polymer gas, provided the amplitudes satisfy the bound [Bry]:

$$\sum_{Y \ni 0} |A(Y)| e^{|Y|} < 1 \quad (\text{III.1.19})$$

The condition  $Y \ni 0$  is there to break translation invariance, and if (III.1.19) holds uniformly in  $\Lambda$ , it has to hold in an infinite volume as well. Before continuing with the Mayer expansion we will show why we can achieve estimates even quite stronger than (III.1.19), using the smallness of the coupling constant, and the fast decay of the propagators. Hence let us prove:

### **Theorem III.1.1: The polymer bound**

Let  $K$  be any fixed constant. For  $g$  small enough we have:

$$\sum_{Y \ni 0} |A(Y)| K^{|Y|} < 1 \quad (\text{III.1.20})$$

where the sum in (III.1.20) is over all finite polymers in an infinite volume (because the estimate (III.1.20) in finite volume  $\Lambda$  is independent of the volume of  $\Lambda$ ).

We explain in detail the proof, restricting for simplicity to vacuum polymers \*. A crucial first step is to go from ordered trees to regular, unordered trees. Let us call  $T(T')$  the regular (unordered) tree associated to  $T'$ . Battle and Federbush realized that at fixed  $T$  there is nothing in (III.1.16) allowing to sum over the various orderings of the trees  $T'$  with  $T(T') = T$ , except the  $s$  dependence in (III.1.17). The following lemma is inspired by their work [BaF1]:

**Lemma III.1.1** For any (regular, unordered) tree  $T$  we have, summarizing all the  $s$  integrals in (III.1.16) by the notation  $\int ds$ :

$$\sum_{T', T(T')=T} \int M_{T'}(s) ds = 1 \quad (\text{III.1.21})$$

This lemma is often stated as the fact that  $\sum_{T', T(T')=T} M_{T'}(s) ds$  is a probability measure  $dp_T(s)$  [Bry].

**Proof** Let us introduce a parameter  $\epsilon_{i,j}$  for each pair of cubes  $\Delta_i, \Delta_j$  in  $Y$ . We can apply the inductive analysis above to the function  $F(\epsilon) = \prod_{(i,j)} (1 + \epsilon_{i,j})$ . More explicitly starting from the root  $i_1 = 1$  we introduce a first interpolating parameter:

$$\epsilon_{i,j}(s_1) = \epsilon_{i,j}[s_1 + (1 - s_1)[\delta_{1,i}\delta_{1,j} + (1 - \delta_{1,i})(1 - \delta_{1,j})]] \quad (\text{III.1.22})$$

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\* For polymers with external variables the proof is similar except the position of the external variables can be used to break translation invariance, hence there is no need to include a condition like  $Y \ni 0$  in (III.1.20).

and perform a first order Taylor expansion  $F(1) = F(0) + \int_0^1 \frac{dF}{ds_1} ds_1$ , and so forth.

The polynomial  $s$  dependence generated by this process is again  $M_{T'}(s)$  for an ordered tree  $T'$ ; but by comparing powers of  $\epsilon_{i,j}$  we can now compute that the left hand side of (III.1.21) is exactly the coefficient of  $\prod_{(i,j) \in T} \epsilon_{i,j}$  in  $F$ , hence is 1.

Remember that isolated cubes are excluded, so if  $Y = \{\Delta_1, \dots, \Delta_p\}$   $p > 1$ , we compute now the action of the  $\prod \left( \frac{\delta}{\delta \phi(x_k)} \frac{\delta}{\delta \phi(y_k)} \right)$  in (III.1.16). This action is a bit complicated to write down, but remember that  $\frac{\delta}{\delta \phi(y_k)}$  acts in  $\Delta_{j_k}$ , which is  $Y_k(T') - Y_{k-1}(T')$  with previous notation; hence it has to derive a vertex from the exponential by the formula

$$\frac{\delta}{\delta \phi(y)} e^{-g \int_{\Delta} \phi^4} = -4g\phi^3(y) e^{-g \int_{\Delta} \phi^4}. \quad (\text{III.1.23})$$

This remark ensures that in the bound for  $A(Y)$  we can extract a factor at least  $g^p$ , and since  $g$  is as small as we want we have a “small factor per cube”; in particular this remark takes care of the large constant  $K$  in Theorem III.1.1. More precisely we claim that:

**Lemma III.1.2** For a propagator decaying like (III.1.6) there exists some large constant  $K' = \text{const.} K$  such that

$$K^p |A(Y)| \leq (K')^p g^p \sum_T \prod_k \left[ \frac{1}{1 + \text{dist}(\Delta_{i_k}, \Delta_{j_k})} \right]^{(3r/4)} \prod_{i=1}^p (d_i!)^{3/2} \quad (\text{III.1.24})$$

where  $d_i$  is the coordination number of the (regular, unordered) tree  $T$  at cube  $\Delta_i$ , i.e., the number of propagators of the tree  $T$  which hook to a vertex within  $\Delta_i$ .

In the case of exponential decay (III.1.5), simply replace  $\prod_k \left[ \frac{1}{1 + \text{dist}(\Delta_{i_k}, \Delta_{j_k})} \right]^{(3r/4)}$  by  $e^{-(3/4) \sum_k \text{dist}(\Delta_{i_k}, \Delta_{j_k})}$ .

Indeed after performing the  $\frac{\delta}{\delta \phi}$  functional derivatives we obtain a sum over procedures  $P$  of functional integrals of the type  $\int \prod_j \phi_j e^{-g \int_Y \phi^4} d\mu_{\{s\}}(\phi)$ ; the product  $\prod_j \phi_j$  is a short notation for a product  $\prod_{\Delta \in Y} \prod_{k=1}^3 [\int_{\Delta} dx \phi^k(x)]^{n_k}$  of integrals over  $\Delta$  of products of one up to three fields which were produced by the  $\frac{\delta}{\delta \phi}$  functional derivatives. We commute the spatial integrals  $dx$  and the functional integral  $d\mu$  and evaluate the latter by a Schwarz inequality:

$$\begin{aligned} \int \prod_j \phi_j e^{-g \int_Y \phi^4} d\mu_{\{s\}} &\leq \left[ \int \prod_j \phi_j^2 d\mu_{\{s\}}(\phi) \right]^{1/2} \left[ \int e^{-2g \int_Y \phi^4} d\mu_{\{s\}}(\phi) \right]^{1/2} \\ &\leq \left[ \int \prod_j \phi_j^2 d\mu_{\{s\}}(\phi) \right]^{1/2} \leq \left[ \int \prod_j \phi_j^2 d\mu(\phi) \right]^{1/2} \end{aligned} \quad (\text{III.1.25})$$

The last inequality is true because both sides can be computed with Wick’s theorem, and  $C(s, x, y)$  is bounded pointwise by  $C(x, y)$ . This is because it is a convex

combination of covariances pointwise bounded by  $C$  (see (III.1.8) and (III.1.12)), and these covariances are pointwise positive by (I.3). Remark that it is a nice coincidence that  $C$  is both pointwise positive and of positive type; in general it is *not* true that a convex combination like (III.1.8) is bounded in the operator sense by  $C$ .

The measure  $d\mu$  has both ultra violet and infra red cutoff and the number  $q_i$  of fields  $\phi_i$  localized in a cube  $\Delta_i$  has to be linear in  $d_i$ , in fact bounded by  $3d_i$  (each derivation step produces at most in fact 3 new fields). There is no longer permutation or symmetry factor between the vertices created in  $\Delta_i$ . At an intermediate stage, after  $j$  steps ( $1 \leq j \leq d_i$ ) of the cluster expansion, the number of fields in  $\Delta_i$  is  $q_{i,j} \leq q_i$ . At step  $j+1$  a new propagator which hooks to  $\Delta_i$  can hook to a field already produced, or it can derive a new vertex. In the first case we have to pay a factor  $q_{i,j} \leq q_i$ , but  $q_{i,j+1} = q_{i,j}$ ; in the second case there is only a factor 1 (coupling constants  $g$  have been already taken into account), but  $q_{i,j+1} = q_{i,j} + 3$ . In the end the gaussian functional integral (III.1.25), by an analogue of the local factorial principle (Lemma II.6.2), gives a factor  $\prod_i (q_i!)^{3/2}$ .

To use Lemma II.6.2 we need a (summable) piece of the spatial decay of the propagator, e.g. one fourth of the initial decay (III.1.6), which explains the factor  $3r/4$  in (III.1.24).

Summing over all possibilities for the functional derivations  $\frac{\delta}{\delta\phi}$ , called also *procedures*  $P$ , the final bound is therefore similar to the worst scenario where each propagator derives a new vertex, hence is  $(\text{const})^n \cdot \prod_i (d_i!)^{3/2}$ . Using Lemma III.1.1 (since the  $s$  dependence is now factorized) we can perform the  $s$  integrals and change the sum over  $T'$  into a sum over regular trees  $T$ . This completes the proof of (III.1.24) (the term  $(Z(\Delta))^{-p}$  is absorbed in the constant  $c$  of (III.1.24)).

Constant powers of factorials like  $(d_i!)^{3/2}$  can be beaten by the decrease of the propagators (in fact with only a piece, say half of it) because of a volume effect: it relies on the fact that the  $d_i$  cubes hooked to  $\Delta_i$  by the tree  $T$  have to be all *distinct*, hence when  $d_i$  gets large, since we are in a finite dimensional space, many of these cubes have to be quite far from  $\Delta_i$ .

**Lemma III.1.3** For any constant  $c$ , taking  $r$  large enough (depending on  $c$ ) we have

$$\prod_{k=1}^{p-1} \left[ \frac{1}{1 + \text{dist}(\Delta_{i_k}, \Delta_{j_k})} \right]^{r/4} \prod_{i=1}^n (d_i!)^c \leq 1 \quad (\text{III.1.26})$$

**Proof** By the volume argument at least half of the  $d_i$  distinct cubes hooked to  $\Delta_i$  have to be at distance at least  $(a \cdot d_i)^{1/d}$ , where  $a$  is a numerical constant and  $d$  is the space time dimension. From the corresponding  $\prod \left( \frac{1}{1 + \text{dist}(\Delta_{i_k}, \Delta_{j_k})} \right)^{r/4}$  we can therefore extract a factor

$$\prod_i \left( \frac{1}{1 + ad_i} \right)^{\frac{r d_i}{16}} \quad (\text{III.1.27})$$

(the factor 16 is because we take only half of the  $d_i$  cubes into account, and a propagator hooks to two cubes). Obviously for  $r$  large enough this leads to (III.1.26).

Theorem III.1.1 follows then easily from one further lemma:

**Lemma III.1.4** For  $r' > d$ , there exists a constant  $K'$  such that:

$$A = \sum_{\substack{Y \ni 0 \\ |Y|=p}} \sum_T \prod_{i=1}^n d_i! \prod_k \left[ \frac{1}{1 + \text{dist}(\Delta_{i_k}, \Delta_{j_k})} \right]^{r'} \leq K'^p. \quad (\text{III.1.28})$$

**Proof** We can interchange the sum over  $Y$  and  $T$ . By translation invariance we can require the root of the tree to be the cube containing the origin. Knowing the structure of  $T$  we can sum over the positions of the cubes of  $Y$ , using the decay in (III.1.28); this would even be true with any summable decay in (III.1.28) (not necessarily power-law with a large power). The result is bounded by  $c^p p!$  because the set  $Y$  gets counted  $p!$  times in the independent summation over its elements. Using Cayley's theorem (section I.4), we can perform the sum over trees  $T$  with given coordination numbers and get:

$$A \leq \sum_{\{d_1, \dots, d_p\}, d_i \geq 1, \sum d_i = 2p-2} c_1^p \leq K'^p \quad (\text{III.1.29})$$

with  $K' = 4c_1$ .

Combining Lemmas III.1.1-4 we have finally to sum in Theorem III.1.1 a series bounded by the geometric series  $\sum_{p>1} (c' K K' g)^p$ , and taking  $g$  small enough achieves the proof of (III.1.20).

Remark that a sloppier version of Lemma III.1.4 without factors  $d_i!$  in (III.1.28) (and without Cayley's theorem) would have suffice at this stage. Lemma III.1.3 in fact takes care of any factorials of the coordination numbers when a large decay is available. However it is a legitimate question to ask whether Theorem III.1.1 remains true under weaker assumptions than (III.1.5-6). The typical rule of thumb is that a cluster expansion usually does not require more than a summable propagator in  $x$ -space. In particular let us sketch the proof of a more powerful theorem:

### Theorem III.1.2: Generalization of Theorem III.1.1

Theorem III.1.1 also holds under the assumption (III.1.6) provided only that  $r > d$  (summable decay).

With such a limited decay, we have no analogue of Lemma III.1.3, and we must be careful not to consume any decay in Lemma III.1.2, since all the decay of the propagator should be kept for the equivalent of Lemma III.1.4. Comparing Lemmas III.1.2 and III.1.4 it seems that we are going to loose the game anyway because there is a power  $3/2$  instead of  $1$  for the local factorials. But there is one point on which we can improve: the Schwartz inequality (III.1.25) is not optimal. We can instead use a Hölder inequality, which we state only in the “worst case”:

$$\int_{\Delta} \phi(x)^3 dx C(x, y) \leq \left[ \int_{\Delta} \phi^4(x) dx \right]^{3/4} \left[ \int_{\Delta} dx C(x, y)^4 \right]^{1/4} \quad (\text{III.1.30})$$

The  $q$  vertices in  $\Delta$  are transformed in this way. The integrated propagators  $[\int_{\Delta} dx C(x, y)^4]^{1/4}$  can be used to sum over  $\Delta$  just as well as  $C(x, y)$  if (III.1.6) holds with  $r > d$ . Hence Lemma III.1.4 remains true. But we can now use the

fact that  $x^n e^{-x} \leq n!$  to improve our factorials: in the worst case there are  $d_i$  vertices to which (III.1.3) is applied in  $\Delta_i$ . They are bounded using the inequality  $g^{d_i/4} \cdot [\int_{\Delta_i} g\phi^4]^{3d_i/4} e^{-g \int_{\Delta_i} \phi^4} \leq c \cdot g^{d_i/4} (d_i!)^{3/4}$ , and (in contrast with Lemma III.1.2) this does not consume any fraction of the propagator's decay, because this bound is completely local (works separately in each cube  $\Delta_i$ ). Comparing to Lemma III.1.4 we win now the game provided  $g$  is taken still smaller than before. We leave to the reader to fill in the details, in particular to check that the sum over all possibilities is of the same order than the “worst case” considered. The idea of using the interaction to improve over gaussian integration is the first example we meet of “domination”, a technique to be discussed at length in the next chapter, hence we do not develop it here in full detail.

Before going on to the Mayer expansion, let us apply the cluster expansion formalism to the proof of Lemma II.6.3, which we postponed until now.

**Proof of Lemma II.6.3** We apply the cluster expansion no longer to  $e^{-g \int_Y \phi^4}$  but to  $(1/n!)(\int_Y \phi^4)^n$ . The number of occupied cubes (vertex domain in the language of section II.6) is still  $p$ . We repeat exactly the same analysis, including the analogue of the Schwarz inequality (III.1.25) to separate the fields hooked to propagators derived by the  $\frac{\delta}{\delta\phi}$  operators from the remainder, which is necessarily of the form  $(1/q!)(\int_Y \phi^4)^q$  for some  $q$  with  $0 \leq q \leq n-p$ . The fields of the first kind are localized in particular cubes of  $Y$  and again their Wick contractions give only factorials of the coordination number, which are bounded like in Lemma III.1.3. Hence the only difference is in the second factor of the first line of (III.1.25), which is no longer bounded by 1, but by  $K^q q!$ . This, together with the obvious inequality  $q! \leq \frac{n!}{p!}$  for  $q \leq n-p$  achieves the proof of Lemma II.6.3.

## B) The Mayer expansion

The Mayer expansion starts with formulas (III.1.15) or (III.1.18), and allows us to compute the correct quantity for a thermodynamical limit, namely normalized Schwinger functions or the pressure. Let us call a finite ordered sequence of polymers such as  $Y_1, \dots, Y_q$  a configuration  $M$ , of length  $q$  (the terminology “Mayer graph” in [FMRS4] is not very appropriate). When there are external variables in  $Y_1, \dots, Y_q$ , the union of which is  $\{z\} = \{z_1, \dots, z_N\}$  we call  $M$  a  $\{z\}$ -configuration. We introduce also the set  $\mathbf{P}(M)$  of all pairs  $(i, j)$ ,  $1 \leq i < j \leq q$ . We say that the configuration is disjoint if  $Y_i \cap Y_j = \emptyset$  for every  $(i, j) \in \mathbf{P}(M)$ . Conversely it is called connected if for every  $(i, j) \in \mathbf{P}(M)$  we can find a chain  $(Y_{i_1} = Y_i, Y_{i_2})(Y_{i_2} Y_{i_3}) \cdots (Y_{i_{k-1}}, Y_{i_k} = Y_j)$  of overlapping polymers in the configuration which join  $Y_i$  to  $Y_j$ , i.e. which are such that  $Y_{i_\ell} \cap Y_{i_{\ell+1}} \neq \emptyset$ . Defining the amplitude of the configuration as  $A(M) = \frac{1}{q!} \prod_{i=1}^q A(Y_i)$  and defining the two body hardcore interaction  $V(Y, Y')$  as 0 if  $Y$  and  $Y'$  are disjoint and  $+\infty$  if they overlap, we can rewrite (III.1.15) or (III.1.18) as

$$\begin{aligned} \tilde{Z}(\Lambda) &= \sum_{\text{disjoint } \emptyset\text{-configurations } M} A(M) \\ &= \sum_{\emptyset\text{-configurations } M} A(M) \prod_{(i,j) \in \mathbf{P}(M)} e^{-V(Y_i, Y_j)} \end{aligned} \quad (\text{III.1.31})$$

$$\begin{aligned} \widetilde{S_{N,\Lambda}^u}(z_1, \dots, z_N) &= \sum_{\text{disjoint } \{z\}-\text{configurations } M} A(M) \\ &= \sum_{\{z\}-\text{configurations } M} A(M) \prod_{i,j \in \mathbf{P}(M)} e^{-V(Y_i, Y_j)} \end{aligned} \quad (\text{III.1.32})$$

(III.1.31) is the grand canonical partition function of a gas of polymers with hard-core interaction: an ideal gas of polymers would have no such interaction. The Mayer expansion is an expansion for the logarithm of (III.1.31) around the ideal gas situation. In our case it consists simply in writing  $e^{-V(Y_i, Y_j)} = (e^{-V(Y_i, Y_j)} - 1) + 1$ . The  $+1$  term corresponds to the free “ideal gas” where the two polymers  $Y_i$  and  $Y_j$  are summed independently (this intuitively restores for  $Y_i$  the “translation invariance” which was broken by the forbidden region  $Y_j$ ); and the  $(e^{-V(Y_i, Y_j)} - 1)$  which is  $-1$  if  $Y_i$  and  $Y_j$  overlap and 0 otherwise is called a “Mayer link” between polymers; it plays indeed the role of a connection somewhat similar to the propagators which link two cubes between a given polymer.

Just as for the cluster expansion we could write a systematic expansion of the “pair of cubes” type, namely write

$$\prod_{\substack{i,j \in \mathbf{P}(M) \\ i < j}} e^{-V(Y_i, Y_j)} = \sum_{J \subset \mathbf{P}(M)} \prod_{(i,j) \in J} (e^{-V(Y_i, Y_j)} - 1) \quad (\text{III.1.33})$$

The result is then factorized over maximal subsequences of  $M$  which are connected by the bonds of  $J$ . Such subsequences can be considered again as configurations which have to be connected. Their connected amplitude is simply:

$$A^T(M) = T(M) \cdot A(M) \quad \text{with} \quad T(M) = \sum_{J \in J^T(M)} \prod_{(i,j) \in J} (e^{-V(Y_i, Y_j)} - 1) \quad (\text{III.1.34})$$

where  $J^T(M)$  is the set of all subsets of  $\mathbf{P}(M)$  connecting  $M$  into a single component. Using simple multinomial identities we obtain that

$$\tilde{Z}(\Lambda) = \sum_{n=0}^{\infty} \frac{1}{n!} \left\{ \sum_{\emptyset-\text{configurations } M} A^T(M) \right\}^n \quad (\text{III.1.35})$$

so that

$$\tilde{p}(\Lambda) = \frac{1}{|\Lambda|} \log \tilde{Z}(\Lambda) = \frac{1}{|\Lambda|} \sum_{\emptyset-\text{configurations } M} A^T(M) \quad (\text{III.1.36})$$

The  $1/n!$  in (III.1.35) comes from the number of partitions of the  $q$  elements of the sequence  $Y_1, \dots, Y_q$  into  $n$  subsequences. In (III.1.34) or (III.1.36) only connected configurations contribute, otherwise  $A^T(M) = 0$ . Convergence of (III.1.36) is not easy to prove in this direct “brute force” expansion because naive estimates which do not take into account the signs cancellations in the factor (III.1.36) fail. (see e.g. [Se]). Again the “pair of cubes” is not the minimal process to reach factorization. This minimal process is rather of the tree type, and we explain it now. Starting from a configuration  $Y_1, \dots, Y_q$ , it is better to first expand  $\prod_{j>1} e^{-V(Y_1, Y_j)}$  in the left hand side of (III.1.33) as:

$$\prod_{j>1} e^{-V(Y_1, Y_j)} = 1 + \sum_{j_1 \neq 1} (e^{-V(Y_1, Y_{j_1})} - 1) \prod_{k>j_1} e^{-V(Y_1, Y_k)}. \quad (\text{III.1.37})$$

As is the rule with a tree expansion the next step depends on which term is selected in the sum (III.1.37). The term 1 frees  $Y_1$  from every hardcore interaction with other polymers, in which case we can directly replace  $Y_1$  by  $Y_2$  and apply an analogue of (III.1.37) to  $\prod_{j>2} e^{-V(Y_2, Y_j)}$ . But for a term corresponding to  $j_1$  in the sum (III.1.37),  $Y_1$  and  $Y_{j_1}$  are linked through a “Mayer link”. The next step is to expand all the remaining constraints relative to  $Y_1$  and  $Y_{j_1}$ , to test whether  $Y_1 \cup Y_{j_1}$  is free from hardcore interaction with the rest of the configuration or not. This is done with priority to  $Y_1$ , i.e. we write

$$\begin{aligned} & \prod_{k>j_1} e^{-V(Y_1, Y_k)} \prod_{\substack{1 < k' \\ k' \neq j_1}} e^{-V(Y_{j_1}, Y_{k'})} = \\ & 1 + \sum_{j_2 > j_1} (e^{-V(Y_1, Y_{j_2})} - 1) \prod_{k>j_2} e^{-V(Y_1, Y_k)} \prod_{\substack{1 < k' \\ k' \neq j_1}} e^{-V(Y_{j_1}, Y_{k'})} \\ & + \sum_{j_2 \neq 1, j_1} (e^{-V(Y_{j_1}, Y_{j_2})} - 1) \prod_{\substack{j_2 < k' \\ k' \neq j_1}} e^{-V(Y_{j_1}, Y_{k'})} \end{aligned} \quad (\text{III.1.38})$$

The first term frees  $Y_1$  and  $Y_{j_1}$  from the rest, the second term corresponds to the tree  $\{(Y_1, Y_{j_1}), (Y_1, Y_{j_2})\}$ , the third one to the tree  $\{(Y_1, Y_{j_1}), (Y_{j_1}, Y_{j_2})\}$ . We can continue this process until all polymers are exhausted, but there are several different possible strategies, among which the strategy “push each branch as far as possible” (also called “turn around the tree”) and the strategy “let the tree grow layer by layer”. The second one is perhaps more natural, so let us adopt it. It means that we develop in priority all the constraints of  $Y_1$  with other polymers, until the process stops by the choice of a factor 1. Then we have built a “first layer” of polymers  $Y_{j_1}, \dots, Y_{j_k}$  linked to  $Y_1$  by Mayer links, with  $j_1 < j_2 < \dots < j_k$ . We expand then in priority the constraints of  $Y_{j_1}$  with all polymers other than  $Y_1, Y_{j_2}, \dots, Y_{j_k}$ , constructing the piece of the “second layer” of the tree linked to  $Y_{j_1}$ . Then we expand in priority the hardcore constraints of  $Y_{j_2}$  with all other remaining polymers, constructing a second piece of the second layer, and continue in this way until all the second layer has been built. Then in this second layer we select the polymer with lowest index in the initial ordering, and expand in priority its constraint with the remaining ones, and continue until all the third layer is built, and so on.

In this way each time the process stops by selecting a term 1 in the sums similar to (III.1.38) for all the polymers of the last layer, we have factorized a particular subsequence. When all polymers are exhausted we end up with independent sums over factorized subsequences connected by Mayer links just as in (III.1.35). Therefore the two expansions must be identical in the sense that they express finally  $\tilde{Z}(\Lambda)$  or  $\tilde{p}(\Lambda)$  by the same series of the same connected amplitudes  $A^T(M)$  for configurations  $M$ ; the advantage of the second point of view lies in the fact that it gives an explicit rewriting of the factor  $T(M)$  as a sum of factors which have absolute values less than 1. More precisely if we study with care (III.1.38) and its generalizations we conclude that it generates exactly once each tree built on the configuration; it does not generate ordered trees as the former

cluster expansion because at a given layer the constraints are always developed in the order provided by the initial ordering  $Y_1, \dots, Y_q$  of the configuration; we see for instance that the tree  $\{(Y_1, Y_{j_1}), (Y_1, Y_{j_2})\}$  is generated once, not twice, because of the condition  $j_2 > j_1$  in (III.1.38). Obviously the total factor associated to a particular tree is in absolute value bounded by 1; this was *not* the case for  $T(M)$  in (III.1.34). However the precise description of this factor is a bit cumbersome. It is the product over all pairs of polymers of the tree of:

- a Mayer link  $e^{-V} - 1$  if the polymers are joined by a line of the tree,
- a hardcore constraint  $e^{-V}$  if the polymers  $Y_i$  and  $Y_j$  belong to the same layer of the tree or if they belong to adjacent layers, e.g. respectively layer  $k$  and  $k + 1$ , and the ancestor  $Y'_j$  at level  $k$  to which  $Y_j$  is hooked has index smaller than  $Y_i$ ,
- a factor 1 otherwise

We conclude that the process we describe ((III.1.37-38) and its generalization) is an algebraic way of reorganizing the sum over  $J^T(M)$  (III.1.34) as a sum over trees, by grouping together many connecting subsets  $J$  having a particular tree in common, according to a particular rule. The details of the rule (“layer by layer”) are partly arbitrary. The gain lies in the fact that although many  $J$ ’s are grouped together, the corresponding large sum of factors  $\prod(e^{-V} - 1)$  is still bounded by 1, so we have effectively taken into account the sign cancellations in (III.1.34).

Having clarified this point, we proceed with the evaluation of the series (III.1.36). We have to bound

$$\sum_{Y_1, \dots, Y_q} \frac{1}{q!} \sum_{T \text{ connecting the } Y_i's} \prod_{j=1}^q A(Y_i) \quad (\text{III.1.39})$$

We can use Cayley’s theorem (Sect.I.4) on the number of tree graphs with fixed incidence numbers  $d_1, \dots, d_q$ : this gives

$$\begin{aligned} \left| \sum_T (\cdot) \right| &= \left| \sum_{d_1, \dots, d_q} \sum_{T, \{d_i\} \text{fixed}} (\cdots) \right| \\ &\leq \sum_{d_1, \dots, d_q} \frac{(q-2)!}{\prod_{j=1}^q (d_i - 1)!} \sup_{(T, d) \text{ fixed}} |(\cdot)| \end{aligned} \quad (\text{III.1.40})$$

Taking into account the  $\frac{1}{q!}$  in (III.1.39), it remains, for a fixed tree connecting the  $Y_i$ ’s, to sum over the  $Y_i$ ’s, not forgetting the important  $\left[ \prod_{i=1}^q (d_i - 1)! \right]^{-1}$  from (III.1.40). We sum over the  $Y_i$ ’s starting from the end branches of the tree (again,  $Y_1$  being the last polymer to be summed, or the “root” of the tree). If  $Y_q$  is such an end branch, let  $\Delta_q$  be a cube in the non empty intersection  $Y_j \cap Y_q$  of  $Y_q$  with its immediate ancestor  $Y_j$  in the tree (remember that  $Y_j$  and  $Y_q$  have to overlap since they are joined by a tree branch, which corresponds to a “Mayer link”). We sum over  $Y_q$ , holding  $\Delta_q$  fixed. This produces, by translation invariance, a factor

$$\sum_{Y \ni 0} |A(Y)| = \sum_{Y \ni 0} A(Y) |Y|^{d_q - 1} \quad (\text{III.1.41})$$

since  $d_q = 1$ . Now we can sum over  $\Delta_q$  in  $Y_j$ , obtaining a factor  $|Y_j|$ . Iterating this process, we can perform inductively the sums over  $Y$ 's, progressively stripping off the branches of the tree. A sum over a given  $Y_i$  gives rise to a factor

$$\sum_{Y \ni 0} |A(Y)| |Y|^{d_i - 1} \quad (\text{III.1.42})$$

where  $|Y|^{d_i - 1}$  arises because at the time  $Y_i$  is summed there has been  $d_i - 1$  other polymers already summed which had  $Y_i$  as their ancestor in the tree. We take into account the crucial factors  $\frac{1}{(d_i - 1)!}$  by writing  $\sum_{d_i} \frac{|Y|^{d_i - 1}}{(d_i - 1)!} = e^{|Y|}$ . Finally in the last sum over  $Y_1$ , the last reference cube  $\Delta_1$  can be anywhere in  $\Lambda$ , hence we get the estimate:

$$\sum_{\emptyset\text{-configurations } M} |A^T(M)| \leq |\Lambda| \cdot \sum_q \left[ \sum_{Y \ni 0} |A(Y)| e^{|Y|} \right] \quad (\text{III.1.43})$$

hence using Theorem III.1.1 we conclude that for  $g$  small enough the series (III.1.36) for the pressure are absolutely convergent.

We can extend easily this analysis to the normalized Schwinger functions. The main difficulty is that the external variables must then be part of the definition of the polymer, so that strictly speaking the polymer is now a set of cubes (its support) plus a set of external variables, (the corresponding amplitude being zero if the support of the smearing functions for the external variables does not intersect the support of the polymer). One difference is that in the cluster expansion generating the polymers some integrations by parts may “hook” to the external fields instead of deriving a vertex by formula (III.1.23); this is no problem for convergence if the number of external variables is bounded, as is the case here. An other subtlety is that one should take into account the fact that polymers which contain external variables are indexed by them, so the corresponding sums have no longer to be symmetrized by  $\frac{1}{n!}$  factors as for vacuum polymers. Therefore:

$$\widetilde{S}_{\Lambda, N}^u \{z_1, \dots, z_N\} = \sum_{\omega = (\omega_1, \dots, \omega_m)} \prod_{i=1}^m \sum_{M \text{ } \omega_i\text{-configuration}} A^T(M) \cdot \sum_{n=0}^{\infty} \frac{1}{n!} \left\{ \sum_{M \text{ } \emptyset\text{-configuration}} A^T(M) \right\}^n \quad (\text{III.1.44})$$

where the sum is taken over partition  $\omega$  of all external variables into  $m$  subsets  $\omega_1, \dots, \omega_m$ . In (III.1.44) we recognize that  $\widetilde{Z}(\Lambda)$  is factorized, so that the normalized functions are given by:

$$S_{\Lambda, N}(z_1, \dots, z_N) = \sum_{\omega = \{\omega_1, \dots, \omega_m\}} \prod_{i=1}^m \sum_{M \text{ } \omega_i\text{-configuration}} A^T(M) \quad (\text{III.1.45})$$

and the normalized *truncated* Schwinger functions are given by the simpler expansion:

$$S_{\Lambda, N}^T(z_1, \dots, z_N) = \sum_{M \text{ } \{z_1, \dots, z_N\}\text{-configuration}} A^T(M) \quad (\text{III.1.46})$$

The series (III.1.45) and (III.1.46) are shown to be absolutely convergent exactly as the one for the pressure, except that now the external variables (which for simplicity are localized by the smearing test functions within given cubes of  $\Lambda$ ) break translation invariance and can be used to hook the last summation over  $Y_1$ : for the pressure this last summation was free, hence the dividing factor  $\frac{1}{|\Lambda|}$  had to be included.

Remark that if we are interested only into the normalized (*not truncated*) Schwinger functions, we may expand only the hardcore constraints between pairs of polymers for which at least one is a vacuum polymer, keeping the constraints between polymers with external variables unexpanded. This will be enough to factorize  $\tilde{Z}(\Lambda)$  as in (III.1.44). We may also derive intermediate versions, in which e. g. the constraints involving at least one vacuum, two or four point polymer are expanded, but not the other. This results in partly truncated amplitudes, which are useful when renormalization of two and four point functions is involved (recall that in part II we learnt that renormalization is best expressed at the level of connected functions). It is a version of this kind which will be used in the next chapters.

The convergence of these various expansions does not lead to any particular problem when the number of external fields is fixed. However in the context of the phase space expansion which we are going to introduce, low momentum fields at a given scale must be considered as external variables. Therefore the number of such variables is no longer bounded at intermediate stages, even if we compute a fixed  $N$ -point Schwinger function. This leads to a subtlety: if we were to expand fully all hardcore constraints involving not only vacuum polymers but also polymers with e.g. two and four external legs, we would generate Mayer configurations in which an arbitrarily large number of external low momentum fields may accumulate at the same place. From the “local factorial principle” we know that this would lead typically to divergent expressions. Therefore the Mayer expansions we use in phase space are slightly modified to avoid this effect: they keep the cubes containing the external legs of a two or four point function (called their external cubes) fixed and non overlapping. In this way the truncation is almost completely performed, and the renormalization cancellations can be performed, but a large number of low momentum fields still cannot accumulate at the same place. This process is described in Section III.3D. Remark that this problem is truly a constructive one, which has no analogue in perturbation theory.

Altogether, we have achieved the proof of:

### **Theorem III.1.3 Convergence of the Mayer expansion**

For  $g$  small enough, the series (III.1.36) and (III.1.45)–(III.1.46) are absolutely convergent, uniformly in  $\Lambda$ . As  $\Lambda \rightarrow \infty$  their limit is therefore still absolutely convergent and can be used to define rigorously the massive weak coupling  $\phi^4$  theory (with ultra violet cutoff) in any dimension.

Let us gather now some further remarks on various aspects of these expansions.

## **C) Further topics**

Our first comment is on the similarities and differences between the cluster and Mayer expansions. In both cases trees emerge as the central structure around

which the expansion is best organized. But for the Mayer expansion there is no positivity requirements to preserve, and the algebraic process (III.1.38) directly generates trees. Ordinary perturbation theory which is nothing but a gas of vertices interacting through propagators, shares these features; there is no positivity requirements, and it is possible to design an “algebraic” or “set-theoretic” cluster expansion based on selecting particular propagators, as sketched in section II.6. But when functional integrals are involved, positivity requirements occur and interpolating parameters seem necessary. The corresponding cluster expansion generate ordered trees rather than trees; however by the Battle-Federbush theorem, the summation from ordered to ordinary trees is controlled by the integration over these interpolating parameters.

An other difference is that in the cluster expansion for the functional integral, strong decay of the propagator and the volume argument allows sloppy estimates on the factorials of coordination numbers of the tree. This is no longer true in the Mayer expansion where one has to use with care the single factorial of this type delivered by Cayley’s Theorem. To get some intuition of this, we may compare the polymers of the cluster expansion to fermions because their hardcore interaction is somewhat reminiscent of the Pauli principle. In contrast, the “Mayer configuration”, made of ordinary polymers with Mayer links, have no longer hardcore interactions; they behave like bosons which can pile up in arbitrary numbers at a given place. Hence we may compare our remark with the known fact that perturbation theory for fermions is more convergent than for bosons, and that keeping track of correct factorial factors for fermions is less important [FMRS4][IM2,Appendix].

Our last comment is on the difference between cluster expansions à la Brydges-Battle-Federbush, used here, and à la Glimm-Jaffe-Spencer as in [GJS]. In section II.6 we wrote the propagator as a Wiener sum over random paths. The difference is summarized by the statement that GJS cluster expansion localizes both vertices and the paths building the propagators, when BBF cluster expansion only localizes the vertices, hence is simpler. However to compute the BBF expansion it is necessary to perform explicitly functional derivatives like (III.1.23). For a polynomial interaction like  $e^{-\int \phi^4}$  this is no problem, but in more complicated situations, it might be difficult or impossible; in such cases one may have to return to GJS cluster expansions (see [Bry] for an example with “large fields holes” where the interaction is non polynomial).

The cluster and Mayer expansions explained above are sufficient for the purpose of proving existence of the thermodynamic limit. They provide information equivalent to connectivity at the graphical level. But sometimes we are interested in more detailed information about the system, in which case one can push further these expansions. A first extension in order to reach information equivalent to graphical one-particle irreducibility was developed in [FMRS4-5] in order to perform full (rather than *useful*) mass renormalization of the models considered there. This point of view has been made more systematic in [IM1-2] for the purpose of multiparticle structure analysis, under the name of p-th order expansions. (Of course the story of the subject is a long one, with some landmarks like [Sp2] or [CFR]). We underline here the basic ideas of the construction, refering the reader to the very clear sections 3 and 4 of [IM1] for a detailed exposition.

The first idea for p-th order cluster expansions is that nothing (at the level

of convergence) prevents us to push the Taylor expansions in the interpolating  $s$  parameters further than to first order. The price to pay however is that the new structures obtained are no longer indexed by ordered trees but by more general “graphs and procedures”. More precisely if we return to interpolation (III.1.8), we start as before with a first order Taylor expansion similar to (III.1.9):

$$Z(\Lambda) = Z(s_1 = 0) + \int_0^1 ds_1 \sum_{\Delta_{1,1} \neq \Delta_1} Z_{\Delta_1, \Delta_{1,1}}(s_1) \quad (\text{III.1.47})$$

In the term  $Z(s_1 = 0)$ ,  $d\mu$  decouples  $\Delta_1$  from its complement  $\Lambda - \Delta_1$ , and as before we turn our attention to the next cube of  $\Lambda - \Delta_1$  and repeat the process. But for the remainder term remark that we call  $\Delta_{1,1}$  the cube (previously called  $\Delta_2$ ) linked by a propagator to  $\Delta_1$ . Since the guiding idea is that a single propagator is not ( $p-1$ )-particle irreducible for  $p \geq 1$ , we are not satisfied and want to push further the expansion, trying to know whether  $\Delta_{1,1}$  is linked to  $\Delta_1$  through more propagators. It is not enough to simply push the Taylor expansion in  $s_1$  up to order  $p$ , because the cubes linked to  $\Delta_1$  by each  $s_1$  derivation may be all different. (This defect would not occur if the cluster was of the pair of cubes type, but we know that this one has a positivity problem). Hence the correct procedure is again very inductive. We introduce a second Taylor formula (hence a new parameter  $s'_1$  interpolating between 0 and  $s_1$ ), writing:

$$Z_{\Delta_1, \Delta_{1,1}}(s_1) = Z_{\Delta_1, \Delta_{1,1}}(0) + \int_0^{s_1} ds'_1 \sum_{\Delta_{1,2} \neq \Delta_1} Z_{\Delta_1, \Delta_{1,1}, \Delta_{1,2}}(s'_1) \quad (\text{III.1.48})$$

In the remainder term  $\Delta_1$  is linked by a propagator both to  $\Delta_{1,1}$  and  $\Delta_{1,2}$ , which of course may coincide. For this remainder term the procedure is pursued until either one variable  $s_1^{(r)}$  is taken at 0, or  $p$  squares among  $\Delta_{1,1}, \dots, \Delta_{1,r}$  coincide. Two typical outcomes of the process are pictured in Fig.III.1.3 (for  $p=4$ ), one in which at most  $p-1$  squares among  $\Delta_{1,1}, \dots, \Delta_{1,r}$  coincide and  $\Delta_1$  is decoupled in the measure  $d\mu$  from  $\Lambda - \Delta_1$ , and the other in which  $\Delta_{1,r} \equiv \Delta_2$  coincides with  $p-1$  previous squares among  $\Delta_{1,1}, \dots, \Delta_{1,r-1}$ . In the first case we choose a new square  $\Delta_2$  in  $\Lambda - \Delta_1$  and continue the expansion, with an important caveat: the new expansion in  $s_2, s'_2, \dots, s_2^{(r_2)}$  is no longer pushed until some  $s_2$  parameter is set to 0 or  $p$  squares linked to  $\Delta_2$  coincide, but until some  $s_2$  parameter is set to 0 or  $\Delta_2$  is linked in a ( $p-1$ )-particle irreducible way to some set of other cubes. This can arrive earlier than for  $\Delta_1$ , because the propagators already created in the first expansion in  $s_1, \dots, s_1^{(r=r_1)}$  may help, as shown in Fig.III.1.4a, which is a possible continuation of the process shown in Fig.III.1.3a.

In the second case (Fig.III.1.3b) we should consider  $\Delta_1$  and  $\Delta_2$  as linked in the same ( $p-1$ )-particle irreducible object and they are treated as a single block in the rest of the expansion, just as were  $\Delta_1$  and  $\Delta_2$  in the ordinary cluster expansion. This means that we introduce an interpolating parameter  $s_2$  by formula (III.1.12). We make first order Taylor expansions in  $s_2, s'_2, \dots, s_2^{(r_2)}$  until either one of these parameters is set to 0 or some square  $\Delta_3$  is linked in a ( $p-1$ )-particle irreducible way to the block  $\Delta_1 \cup \Delta_2$ , as in Fig.III.1.4b.

This process is continued and larger and large ( $p-1$ )-particle irreducible blocks are formed, until the process ends up, an issue guaranteed by the finiteness of  $\Lambda$ .

The result is an expansion in terms of sets of disjoint cubes  $Y_1, \dots, Y_q$  or polymers (again it is better to factorize the functional integral for an empty cube, passing to  $\tilde{Z}$  as in (III.1.14)). For each polymer there is a sum over graphs  $G$  made of lines connecting these cubes. These graphs are no longer trees, but there are some restrictions on these graphs, because they never connect subsets of cubes in a p-particle irreducible way. Finally there is also a sum over procedures  $P$  leading to these graphs, which is a generalization of the sum over orderings of the trees in the ordinary cluster expansion. The sum over procedures leading to the same graph is again controlled by the integration over interpolation parameters, using a generalization of the Battle-Federbush result [IM1, Lemma 1]. The sum over graphs built on the support of a polymer may also seem more difficult than the sum over trees, but the condition that there is no p-particle irreducible structures is in fact very restrictive, so that Theorem III.1.1 remains valid. Therefore p-th order cluster expansion followed by an ordinary Mayer expansion can be used for the computation of thermodynamic limits and normalized quantities. But of course their real interest lies in the fact that in these expansions the k-particle analysis becomes possible. For instance an expansion with  $p=3$  makes explicitly visible all the chains of one-particle irreducible two-point subgraphs in a polymer (this is not the case for  $p=2$ , see Fig.III.1.5).

This opens the possibility of performing full non-perturbative mass renormalization in  $\phi_4^4$ -like models, i.e. to set the renormalized mass of the theory directly at the desired value. Indeed in this renormalization, the one particle irreducible two point functions must be subtracted at 0 external momenta. For the same reason than in perturbation theory, the “useless” part of mass renormalization is not the source of any non-summable effect of the renormalon type. However just as renormalization of the 4 point function requires at least the computation of connected four point structures, hence an ordinary Mayer expansion, the renormalization of such one particle irreducible subgraphs requires a generalized version of the Mayer expansion. Indeed to remove hardcore constraints between connected polymers would not be enough, because inside a connected graph a one particle irreducible two point subgraph has still hardcore constraints with the rest of the graph (hence its 0-momentum value is not “background independent”, but depends on the rest of the expansion, so that it cannot cancel exactly with a universal, background-independent counterterm). The necessary generalized Mayer expansion removes not only the hardcore constraints between connected polymers, but also the hardcore constraints of each one particle irreducible 2-point subgraph inside each connected polymer with the rest of the polymer, by algebraic formulas similar to (III.1.37-38). Here again, to display convergence one has to be careful and to proceed in an inductive way, in the natural order provided by the maximal chains of such one particle irreducible subgraphs [FMRS5]. Although appealing at the conceptual level, this method leads to some rather intricate technicalities, so we decide in section III.3 not to use it and to perform a more standard “fixed point” computation of the renormalized mass.

One can continue along these lines and define still more general “Mayer expansions”. For instance Iagolnitzer and Magnen have considered, in the analysis of the Bethe-Salpeter equation, the removal of hardcore constraints between two-particle irreducible kernels along the chain pictured in Fig.III.1.6 [IM4]. Here again the

linear character of the chain is used to organize the expansion. To our knowledge a completely general theory of such expansions remains to write.

As shown in [IM1–4], p-th order expansions can be used for a better understanding of multiparticle structure and may be of help for a proof of asymptotic completeness which is of course one of the most important open mathematical problems in field theory. For a review on this subject we refer to [Ia].

#### D) The $\phi_2^4$ theory

A single slice model like the one considered above may seem somewhat artificial. We could without difficulty remove the ultraviolet cutoff for the  $\phi_1^4$  theory (the anharmonic oscillator) and apply the previous formalism to construct its weak coupling infinite volume limit. This is still not very exciting. Hence we prefer to conclude this section by a more interesting exercise which gives a first flavor of the ultraviolet problem in constructive field theory. The construction of the  $P(\phi)_2$  model by Nelson, Glimm-Jaffe and Guerra-Rosen-Simon [Er1] was the birth act of constructive theory, hence no book on this subject should omit it. Here we will limit ourselves to rephrase Nelson's probabilistic argument in our language. Using the powerful method of part II for graphical estimates, it is almost immediate to prove the existence of the ultraviolet limit of the theory in a unit cube. The only divergent graphs in the  $\phi_2^4$  theory are due to loops  $\int \frac{d^2 p}{p^2+m^2}$  made of a single line, also called tadpoles. The Wick ordering, or normal ordering of a polynomial with respect to a gaussian measure is precisely the mathematical operation which removes such possibilities, so that the perturbative expansion of Wick ordered quantities does not contain any tadpoles [Si]. Introducing our favorite cutoff (II.1.2-4) for the gaussian measure  $d\mu_\rho$  we expect convergence as  $\rho \rightarrow \infty$  of the  $:\phi^4:$  theory in two dimensions. Let us prove convergence of the normalization:

$$Z_\rho = \int d\mu_\rho(\phi) e^{-g \int_{\Delta} :\phi^4:} \quad (\text{III.1.49})$$

$\Delta$  is a unit square in  $\mathbb{R}^2$ ; the boundary conditions are not very important. The problem is that  $:\phi^4:$  is no longer positive. A simple computation gives:

$$:\phi^4: = \phi^4 - c \cdot \rho \phi^2 + d \cdot \rho^2 \quad (\text{III.1.50})$$

where  $c$  and  $d$  are two numerical constants.  $c\phi^2$  corresponds to the mass renormalization (of the tadpole graph) and  $d$  to the vacuum energy renormalization (the graph with a single vertex and two tadpoles);  $\rho$  is up to a constant the logarithm of the cutoff, since the tadpole diverges logarithmically.

The naive lower bound on (III.1.50) is  $-K \cdot g \cdot \rho^2$  for some constant  $K$ , and the naive bound on  $Z_\rho$  therefore explodes as  $e^{K \cdot \rho^2}$ ; it does not allow to prove the existence of the ultraviolet limit. But the probabilistic measure of the set of fields where this lower bound is reached is in fact very small and therefore  $Z_\rho$  converges [Er1].

Let us prove this using our standard splitting of the propagator as  $C = \sum C^i$ ; we write  $:\phi^4: = :\phi_\rho^4: = \sum_i :\phi^{i,4}:$ , with  $:\phi^{i,4}: \equiv :\phi_i^4: - :\phi_{i-1}^4::$ . To compute  $Z_\rho - Z_{\rho-1}$ , we apply, for each  $i$ ,  $i = 0, \dots, \rho$ , a Taylor expansion in  $:\phi^{i,4}:$  up to

order  $n_i = 1 + a \cdot i$ , where  $a$  is some large integer to be fixed later\*. This means that we write:

$$Z_\rho = \int d\mu_\rho(\phi_\rho) e^{-g \int_{\Delta} \sum_{i=0}^{\rho} t_i : \phi^{i,4}:} |_{\{t=1\}} = \prod_i (I^{n_i} + R^{n_i}) Z_\rho \quad (\text{III.1.51})$$

where the operator  $I^{n_i}$  takes the beginning of the Taylor expansion in  $t_i$  of  $Z$  up to order  $n_i - 1$  and  $R^{n_i}$  takes the Taylor remainder at order  $n_i$ .

In fact we stop expanding the product in (III.1.51) as soon as a remainder term is produced. This means that we write

$$\prod_i (I^{n_i} + R^{n_i}) = \prod_i I^{n_i} + \sum_{j=0}^{\rho} R^{n_j} \prod_{i>j} (I^{n_i}). \quad (\text{III.1.52})$$

We obtain a sum of  $\rho + 2$  terms. In the term with index  $j$  in the sum, the exponential of the interaction has cutoff  $j$ , so it satisfies the bound:

$$e^{-g \int_{\Delta} : \phi_j^4:} \leq K_j, \quad K_j \equiv e^{K \cdot g \cdot j^2} \quad (\text{III.1.53})$$

In the first term  $\prod_i I^{n_i}$  there is no interaction at all; by convention we can consider that it corresponds to  $j = -1$ , with  $R_{n_{-1}} \equiv 1$ , and  $K_{-1} \equiv 1$ .

For  $i > j$  the derivations in the  $I^{n_i}$  operators have produced  $k_i < n_i$  Wick-ordered vertices with highest leg at scale  $i$ , and for  $j \neq -1$  the  $R^{n_j}$  operator produced  $n_j = 1 + a \cdot j$  vertices with highest leg at scale  $j$ . We integrate all these vertices with respect to the gaussian measure, and obtain a sum of graphs with no tadpoles since the vertices are Wick-ordered. These graphs are then bounded with the universal method of part II. In section II.1 we remarked that for the convergent graphs of a superrenormalizable theory (which is the case considered here, since the divergent tadpoles are forbidden), we obtain some vertical exponential decay in index space between the highest scale of each vertex and the 0-th scale. This is different from the just renormalizable convergent case, for which vertical exponential decay occurs only between the highest and the lowest scale of each vertex.

Hence there is a small constant  $\epsilon$  and a large constant  $C_1$  such that if  $n = n_j + \sum_{i>j} k_i = n_j + k$ ,  $k \equiv \sum_{i>j} k_i$  is the total of all vertices produced:

$$|Z_\rho| \leq \sum_{j=-1}^{\rho} K_j \cdot L_j \sum_{\{k_i < 1+a \cdot i\}} (g \cdot C_1)^n n! \prod_{i>j} M^{-\epsilon \cdot i \cdot k_i} \quad (\text{III.1.54})$$

where we recall that  $K_j = e^{K \cdot g \cdot j^2}$  if  $j \neq -1$ ,  $K_{-1} = 1$ . and we define similarly  $L_j \equiv M^{-\epsilon \cdot a \cdot j^2}$  for  $j \neq -1$  and  $L_{-1} \equiv 1$ . The factor  $L_j$  corresponds to the power counting earned for the  $n_j = 1 + aj \geq aj$  vertices produced in slice  $j$ , and the factors  $M^{-\epsilon \cdot i \cdot k_i}$  to the power counting earned for the vertices produced at slices  $i > j$ . This sum is bounded by a constant as  $\rho \rightarrow \infty$ , if  $a$  is taken large enough.

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\* The factor 1 is there so that even in the slice  $i = 0$  an expansion is performed.

Indeed let us choose first  $a$  large so that  $gK < \epsilon \cdot a/2$ . In this case  $K_j \cdot L_j \leq \sqrt{L_j}$ . We use the binomial bound  $n! \leq 2^{k+1+a} j! (1+a)^j$  and define

$$C = 1 + \sum_{j=0}^{\infty} (1+a)^j (2C_1 \cdot g)^{1+a} e^{-\epsilon(a/2)j^2} \leq 1 + |g| C' \quad (\text{III.1.55})$$

to get the bound:

$$|Z_\rho| \leq C \cdot \sum_k \sum_{\{k_i < a \cdot i, \sum k_i = k\}} \frac{(2gC_1)^k}{k!} \left( \frac{k!}{\prod_i k_i!} \right)^2 \prod_i ((k_i!)^2 M^{-\epsilon \cdot i \cdot k_i}) \quad (\text{III.1.56})$$

Since  $k_i < a \cdot i$ , there exists a constant  $C_2$  such that  $((k_i!)^2 M^{-\epsilon \cdot i \cdot k_i}) \leq C_2^{k_i} \cdot M^{-(\epsilon/2)i \cdot k_i}$ . Hence:

$$|Z_\rho| \leq C \cdot \sum_k \sum_{\{k_i, \sum k_i = k\}} \frac{(2gC_1 C_2)^k}{k!} \left( \frac{k!}{\prod_i k_i!} \right)^2 \prod_i M^{-(\epsilon/2)i \cdot k_i} \quad (\text{III.1.57})$$

Using the multinomial identity, we have, if we define  $C_3 = \sum_{i=0}^{\infty} M^{-(\epsilon/4)i}$ :

$$|Z_\rho| \leq C \cdot \sum_k \frac{(2gC_1 C_2 C_3^2)^k}{k!} \leq C \cdot e^{2gC_1 C_2 C_3^2} \quad (\text{III.1.58})$$

hence we have proved a uniform upper bound on  $Z_\rho$  as  $\rho \rightarrow \infty$ . A constant lower bound on  $Z_\rho$  in this particular case is provided in a cheap way by Jensen's inequality:

$$\int d\mu(\phi) e^{-g:\phi^4:} \geq e^{-g \int d\mu(\phi):\phi^4:} = 1 \quad (\text{III.1.59})$$

since the integral of a Wick ordered vertex vanishes.

In fact, applying the same expansion to  $Z_\rho - Z_{\rho-1}$  we would prove in the same way that  $Z_\rho$  is not only bounded, but converges as  $\rho \rightarrow \infty$ . Furthermore by (III.1.55)  $C \leq 1 + |g| C'$ , with  $C'$  a constant, hence the limit  $Z(g) = \lim_{\rho \rightarrow \infty} Z_\rho$  is  $1 + O(g)$  at small  $g$  (in fact  $O(g^2)$ ) would be easy to obtain, again because the integral of a Wick-ordered vertex vanishes).

Altogether we have proved:

#### Theorem III.1.4 Ultraviolet limit of finite volume $\phi_2^4$

For any coupling constant  $g$  we have  $\lim_{\rho \rightarrow \infty} Z_\rho = Z(g) = 1 + O(g)$ , so that  $Z(g)$  is continuous in  $g$  at  $g = 0$  (it is in fact Borel summable [EMS]).

Starting from this result, we may apply the cluster and Mayer expansion to construct the thermodynamic limit of the theory at weak coupling just as in the single slice case discussed above. This is left to the reader.

Of course these results extend without effort to the  $P(\phi)_2$  model for  $P$  any polynomial bounded below. The cluster and Mayer expansions again require small coefficients in this polynomial (which means a high temperature regime in statistical mechanics).

The lesson to be learnt is that to control this first example of a non-trivial ultraviolet limit, one has to push perturbation theory farther and farther at higher

and higher scales. This results in a net gain because the typical integration region of a local object like a vertex is smaller and smaller at higher energy. This idea will be developped in a more optimal way in the multiscale expansion of the next section, in which we introduce a lattice of cubes adapted to each scale of momenta and test for the presence of interaction vertices of the corresponding scale in each cube of this lattice.

### III.2 The phase space expansion: the convergent case.

*The author feels that this technique of deliberate lying will actually make it easier for you to learn the ideas.*

*Once you understand a simple but false rule, it will not be hard to supplement that rule with its exceptions.*

– D. E. Knuth, The *TEXbook*.

#### A. The vertical expansion and convergent polymers

We turn to a description of the natural generalization of the cluster expansion of last section in the case of a model in which the principle of phase space chopping becomes necessary.

What are the models of this type? Consider first the massive  $\phi^4$  theory in dimension  $d$ . For  $d = 2$  the renormalization problem reduces to Wick ordering. We have to do some momentum analysis but we can avoid spatial localization, hence a true phase space expansion as shown in the preceding section. But for  $d = 3$  phase space chopping starts being truly useful [GJ1]. Three dimensional theories are therefore a classical testing ground for the phase space expansionist [Ba1]. It would be interesting for the reader to test also the version of the phase space expansion defined below to  $\phi_3^4$ , and compare it to the original constructions [FO][MS1]. However our formalism is really designed for marginal, just renormalizable theories, so we focus directly on this case. For  $\phi_4^4$ , which is not asymptotically free, the phase space expansion simply proves that if one starts with a bare theory with a bare coupling small enough so that the first cluster expansion at the bare scale converges, then the resulting renormalized theory is a free field. This phenomenon is discussed in the next section (Theorem III.3.2).

Therefore in order to get some non-empty constructive results we have to look for other models: in the next section III.3 we discuss how to apply the phase space expansion to the infra red (critical)  $\phi_4^4$  (with fixed ultraviolet cutoff) and in section III.4 to the massive Gross-Neveu model in two dimensions, which is asymptotically free. In the concluding section III.5 we discuss the much less advanced case of the ultra violet limit of non-abelian gauge theories, where the hope is to gain a non-perturbative understanding of asymptotic freedom.

The phase space expansion is a rather natural extension of the cluster expansion, but it leads to several additional technicalities. Sticking to our principle of dividing the difficulties into pieces easier to digest, we propose in this section an overview of the formalism and a constructive analogue of section II.1, namely an analysis of the convergent cases (with favorable power counting). For this limited purpose our standard model, massive ultra violet  $\phi_4^4$ , is still perfectly convenient.

We want to find the analogues of the phase space slicing of part II, but in a context no longer limited to Feynman graphs; we have to perform true functional integrals. Hence it is not enough to consider propagators only, and we must use the language of fields. We slice the covariance of our gaussian measure, e.g. with  $\alpha$ -space cutoffs as in (II.1.2-4). To this slicing of the covariance is associated the corresponding orthogonal decomposition (II.1.5) of the gaussian measure  $d\mu = \otimes d\mu^i$ . The field  $\phi$  becomes a sum of random variables  $\phi^i$  independently

distributed according to  $d\mu^i$ .  $\phi^i$  will be called a field of frequency (or index)  $i$ , which means that it corresponds to momenta of order  $M^i$ . In each momentum slice  $i$  the gaussian measure  $d\mu^i$  is factorized and obviously one should perform one corresponding cluster expansion, this time with respect to the scaled lattice  $\mathbf{D}^i$  of cubes of side size  $M^{-i}$ . These are called the horizontal cluster expansions. However this is not enough, as in the single scale model, to get the factorizations necessary for the thermodynamic limit, because now the interaction still couples the various horizontal slices. Remember that  $\phi^4$  vertices were pictured as dotted vertical links, dual to the horizontal propagators in Fig.II.1.2. One should find therefore a kind of vertical cluster expansion which is the analogue of the horizontal ones, but for these dual vertical links. Finally there is the question of the Mayer expansion, which we discuss briefly, anticipating on the following.

The simplest solution is to perform first all the horizontal and vertical cluster expansions, hence to obtain a dilute gas of polymers with hardcore interactions swimming in the “ $d + 1$ ” dimensional phase space, and apply a single Mayer expansion to it. This point of view is all right as far as the bare theory is concerned (or for superrenormalizable theories like  $\phi_2^4$  or  $\phi_3^4$ , where only simple mass renormalizations have to be performed). But in the case of just renormalizable theories we want to compute the analogue of the usefully renormalized expansion with effective couplings, whose advantages were detailed at length in Part II. This means that at each scale  $i$  we want to add and subtract a counterterm which in analogy with section II.4 is (in the case of the coupling constant renormalization) a sum over all polymers with 4 external variables made with the fields and propagators of higher slices. To be of the desired form  $\int_{\Lambda} \phi_i^4$  (with  $\phi_i = \sum_{j=0}^i \phi^j$ ), such a counterterm has to live in all of  $\Lambda$  without hardcore constraints. In order to combine this counterterm with the four point polymers, the hardcore constraints between these polymers have to be removed. As shown in the preceding section, this is exactly what the Mayer expansion does. For this reason, we cannot wait for a single Mayer expansion only in the end, and in each new slice after the corresponding horizontal and vertical cluster expansions have been performed, a Mayer expansion has also to be performed.

The drawback is that it is difficult to visualize the result of the corresponding sequence of Mayer expansions. Indeed Mayer configurations are sequences of overlapping polymers, which can at best be pictured as superposed strata. Iterating the superposition of such strata becomes really tedious to picture. This is one of the reasons for which we decide to postpone the corresponding formalism and the problems related to renormalization to the next section, where they will be treated in the concrete case of infrared  $\phi_4^4$ . The goal of this section is to familiarize ourselves with the phase space expansion in the limited context of *convergent* polymers; in this case the sandwiching of Mayer expansions between each cluster expansion is not necessary, and will not be introduced. Hence this chapter is really the constructive generalization of chapter II.1. We will discover the basic mechanism of convergence of the phase space expansion in a context free of the technicalities associated to renormalization; we hope that the experience gained in this way will be valuable to throw light on the general case.

Let us return to the “vertical” cluster expansion. What kind of interpolation

formula should be applied to the vertex  $\int (\sum_i \phi^i)^4$  which mimicks the interpolations of the preceding sections for the propagator? Again the key requirement is that this interpolation preserves positivity. As usual it is convenient to start from the bare theory, hence from the highest momentum slice, and proceed downwards towards the renormalized theory. Therefore the first interpolation should separate  $\phi^\rho$  from the rest, which is  $\sum_0^{\rho-1} \phi^i = \phi_{\rho-1}$ . Also there should be a different interpolation parameter, called  $t_\Delta$  (not to be confused with the  $s$  parameters of horizontal cluster expansions) for every cube  $\Delta$  of  $\mathbf{D}^\rho$ , since these are the units which the  $\rho$ -th horizontal cluster expansion tries to decouple. One of the simplest interpolation satisfying positivity is

$$\int_{\Delta} g(\phi^\rho + \phi_{\rho-1})^4 = [\int_{\Delta} g(\phi^\rho + t_\Delta \phi_{\rho-1})^4 + (1 - t_\Delta^4) \int_{\Delta} g\phi_{\rho-1}^4] |_{t_\Delta=1} \quad (\text{III.2.1})$$

with  $g$  the bare coupling constant. This is an analogue of formula (III.1.8), taking into account that we have now a quartic object (other formulas are of course possible). The  $\rho$ -th vertical expansion consists in applying, after the  $\rho$ -th horizontal cluster expansion the operator:

$$\prod_{\Delta \in \mathbf{D}^\rho} (I + R) \left[ e^{-\sum_{\Delta} \int_{\Delta} g(\phi^\rho + t_\Delta \phi_{\rho-1})^4 + (1 - t_\Delta^4) \int_{\Delta} g\phi_{\rho-1}^4} \right] \quad (\text{III.2.2})$$

where  $I f$  takes the beginning of the Taylor expansion of  $f$  in  $t_\Delta$  up to some fixed order  $p$ , and  $R f$  takes a Taylor remainder  $\int_0^1 dt_\Delta \frac{(1-t_\Delta)^p}{p!} \frac{d^{p+1}}{dt_\Delta^{p+1}} f$ . We allow for the possibility to have  $p > 1$ ; in fact for horizontal cluster expansions a single Taylor step always gives sufficient decay because the sliced propagator has fast decay, but this is not true for the vertical expansion. We have seen that, in the case of  $\phi_4^4$ , it is only for graphs with more than 4 external low momentum legs that exponential decay in the vertical direction is guaranteed. Since each  $\frac{d}{dt_\Delta}$  derivation generates at least one external low momentum leg  $\phi_{\rho-1}$ , in dimension 4 it is wise to choose  $p = 4$ , so that we are sure that the remainder terms correspond to a situation in which vertical exponential decay occurs.

The inductive generalization of (III.2.2) to all scales is easy. Writing  $\mathbf{D}$  for the union of all scaled lattices  $\mathbf{D}^i$  covering  $\Lambda$  with cubes of side size  $M^{-i}$ ,  $0 \leq i \leq \rho$ , we introduce a parameter  $t_\Delta$  for each  $\Delta \in \mathbf{D}$ . Then for fixed  $x \in \Lambda$  we call  $t_i$ ,  $0 \leq i \leq \rho$  the parameter of the cube  $\Delta \in \mathbf{D}^i$  to which  $x$  belongs and we write:

$$\phi_i(t) \equiv \sum_{j=0}^i [ \prod_{j < k \leq i} t_k ] \phi^j \quad (\text{III.2.3})$$

$$\begin{aligned} \phi^4(x) &= \left( \sum_{j=0}^{\rho} \phi^j(x) \right)^4 = \sum_{i=0}^{\rho} (1 - t_{i+1}^4) (\phi_i(t))^4 \Big|_{\{t\}=1} \\ &= \sum_{i=0}^{\rho} (1 - t_{i+1}^4) \left( \sum_{j=0}^i \left[ \prod_{j < k \leq i} t_k \right] \phi^j \right)^4 \Big|_{\{t\}=1} \end{aligned} \quad (\text{III.2.4})$$

where by definition  $t_{\rho+1} \equiv 0$ .

In some cases it may be convenient to use smooth characteristic functions  $\Delta(x)$  for the cubes  $\Delta$ ; then one defines:

$$a(x, \{t\}, i, j) = \prod_{j < k \leq i} \left[ \sum_{\Delta \in \mathbf{D}^k} t_\Delta \Delta(x) \right]; \quad \phi_i(\{t\}, x) = \sum_{j \leq i} a(x, \{t\}, i, j) \phi^j(x) \quad (\text{III.2.5})$$

and one writes

$$\int_{\Lambda} \left( \sum_{i=0}^{\rho} \phi^i(x) \right)^4 dx = \int_{\Lambda} dx \sum_{i=0}^{\rho} \left[ 1 - a(x, t, i, i+1)^4 \right] \phi_i(\{t\}, x) \Big|_{\{t\}=1} \quad (\text{III.2.6})$$

where again, by convention  $t_\Delta \equiv 0$  if  $\Delta \in \mathbf{D}^{\rho+1}$ . Formulas (III.2.5) and (III.2.6) are rather complicated, and we shall not use them in this book; but the use of sharp characteristic functions has some drawback too, and the reader will see below that it forces us to treat a small piece of the gaussian measure in the form of an interaction  $e^{-\epsilon \int \partial_\mu \phi \partial^\mu \phi}$ , with  $\epsilon$  a small constant.

For this reason, and for the treatment of models with counterterms, it is convenient to give also rules to interpolate fields with derivatives and quadratic interactions, like  $m^2 \int \phi^2$  and  $\epsilon \int \partial_\mu \phi \partial^\mu \phi$ . We use the interpolations:

$$\partial_\mu \phi_i(t) \equiv \sum_{j=0}^i \left[ \prod_{j < k \leq i} t_k \right] \partial_\mu \phi^j \quad (\text{III.2.7})$$

$$\phi^2 = \sum_{i=0}^{\rho} (1 - t_{i+1}^2) (\phi_i(t))^2 \Big|_{\{t=1\}} \quad (\text{III.2.8})$$

$$\partial_\mu \phi \partial^\mu \phi = \sum_{i=0}^{\rho} (1 - t_{i+1}^2) [\partial_\mu \phi_i(t) \partial^\mu \phi_i(t)] \Big|_{\{t=1\}} \quad (\text{III.2.9})$$

We should think to these formulas are the inductive generalizations of natural interpolation rules analogous to (III.2.1):

$$\int_{\Delta} m^2 (\phi^\rho + \phi_{\rho-1})^2 = \left[ \int_{\Delta} (m^2) (\phi^\rho + t_\Delta \phi_{\rho-1})^2 + (1 - t_\Delta^2) \int_{\Delta} (m^2) (\phi_{\rho-1})^2 \right] |_{t_\Delta=1} \quad (\text{III.2.10})$$

$$\begin{aligned} \int_{\Delta} \epsilon (\partial_\mu \phi^\rho + \partial_\mu \phi_{\rho-1}) (\partial^\mu \phi^\rho + \partial^\mu \phi^{\rho-1}) &= \left[ \int_{\Delta} \epsilon (\partial_\mu \phi^\rho + t_\Delta \partial_\mu \phi_{\rho-1}) (\partial^\mu \phi^\rho + t_\Delta \partial^\mu \phi_{\rho-1}) \right. \\ &\quad \left. + (1 - t_\Delta^2) \int_{\Delta} \epsilon (\partial_\mu \phi_{\rho-1} \partial^\mu \phi_{\rho-1}) \right] |_{t_\Delta=1} \end{aligned} \quad (\text{III.2.11})$$

Until the expansion at scale  $i$  we should put every parameter  $t_\Delta$  with  $\Delta \in \mathbf{D}_j$ ,  $j \leq i$  to 1; then after the cluster expansion of scale  $i$ , the interpolating parameters  $t_\Delta$  with  $\Delta \in \mathbf{D}_i$  are introduced for all fields (both in the exponential of the interaction and in fields already derived by former expansions). Remark that derivatives and  $t$ -dependence do not commute; indeed the  $t_k$ 's are functions of  $x$ , by our convention that  $t_k$  is the parameter  $t_\Delta$  such that  $x \in \Delta \in \mathbf{D}^k$ . To ensure that  $\partial_\mu \phi^j$  is of order  $M^j \phi^j$ , as it should, it is essential to write the  $t$  dependence of derived fields always after performing the derivative, otherwise the derivative

could act on a characteristic function of a cube of much smaller size, hence be much bigger than what is expected.

The phase space expansion consists in applying at each scale, starting from  $\rho$ , a horizontal cluster expansion, for instance of the tree type as in the preceding section, then the vertical expansion, which means that we apply to the functional integral the operator

$$\prod_{\Delta \in \mathbf{D}^i} \left( I_{\Delta}^{(4)} + R_{\Delta}^{(4)} \right) \quad (\text{III.2.12})$$

where  $I_{\Delta}^{(4)}$  takes the beginning of the Taylor expansion  $t_{\Delta}$  at  $t_{\Delta} = 0$  up to 4<sup>th</sup> order and  $R_{\Delta}^{(4)}$  is the Taylor remainder. To compute the result of this vertical cluster expansion at scale  $i$ , we expand the product (III.2.12), and when the  $I_{\Delta}^{(4)}$  term is chosen we draw a thick line at the bottom of the cube  $\Delta$  as in Fig.III.2.1. This thick line is a sort of vertical analogue of a Dirichlet condition; it suggests the corresponding decoupling of frequencies. We call it also a “closed gate”. When  $R_{\Delta}$  is chosen the corresponding “gate” is open. The cube  $\Delta$  should be thought of as connected through this open gate to the larger cube  $\Delta' \in \mathbf{D}^{i-1}$ ,  $\Delta' \supset \Delta$ . Remark that dotted vertices then cross this open gate, with a total of at least 5 low momentum fields  $\phi_{i-1}$  hooked to them. Such an open gate is also called a *strong* connection, and guarantees automatically a favorable power counting.

Let us make the notion of connectedness in phase space more precise. We say that two cubes  $\Delta \in \mathbf{D}^i$  and  $\Delta' \in \mathbf{D}^j$  are directly connected if either (we may assume  $j \leq i$ ):

- there is a propagator (horizontal line) between them. This is as before. It requires  $i = j$  and the propagator is generated by the  $i$ -th horizontal cluster expansion.
- $\Delta \subset \Delta'$  and there is a vertex (dashed vertical line) localized in  $\Delta$  with two fields hooked to it, one of scale  $i$  and the other of scale  $j$  \*.
- there is an open gate between them. This requires  $i = j + 1$  and  $\Delta \subset \Delta'$ . This open gate has to be generated by the  $i$ -th vertical cluster expansion. The cubes  $\Delta$  and  $\Delta'$  are said to be directly strongly connected.

The first and second conditions are easy generalizations of the horizontal and vertical connections of perturbation theory, but the third one corresponds to a remainder term which has no equivalent in perturbation theory. As remarked, in this third case there are at least 5 low momentum fields localized in  $\Delta$ .

The notion of connectedness is then extended so that cubes  $\Delta_i$  and  $\Delta_j$  which can be joined by a chain of directly connected cubes are connected. Maximal sets

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\* In [FMRS5] an other rule is used: vertical connections which connect  $\Delta$  to  $\Delta'$  are also required to connect together all the cubes  $\Delta''$  which are between  $\Delta$  and  $\Delta'$  (i.e. such that  $\Delta \subset \Delta'' \subset \Delta'$ ; the horizontal analogue would be to require that propagators between two cubes connect together all cubes on a straight line between the two ends of the propagator. Both points of view are perfectly valid and lead to slightly different phase space expansions but which are both convergent (because by power counting the vertical links decay exponentially). But the point of view used here is certainly more natural.

of connected cubes in phase space are called polymers. It is also convenient to define strongly connected domains. Two cubes are said to be strongly connected if they can be linked through a chain of directly strongly connected cubes, and the corresponding maximal sets are the strongly connected domains, pictured in Fig.III.2.2. Hence they are defined exactly like the polymers, but only connections of the third type are taken into account. One should think to these strongly connected domains as to elementary building blocks of the phase space expansion in which the interaction is not decoupled.

To understand better the whole expansion, it is a good exercise to visualize how it works in the first slices.

In the  $\rho$ -slice we perform the first horizontal cluster expansion with the gaussian measure  $d\mu^\rho$  corresponding to fields  $\phi^\rho$  and covariance  $C^\rho$  and with respect to the cubes of  $\mathbf{D}^\rho$ ; then we perform the first vertical expansion. We obtain a sum over a gas of disjoint  $\rho$ -polymers which is an analogue of formula (III.1.15); these polymers are defined again as maximal connected objects. By convention again, it is convenient not to consider the empty isolated cubes as polymers, which leads to a new trivial normalization of the polymers similar to (III.1.14-16). At this stage the only connections to take into account are the propagators of the  $\rho$ -th horizontal cluster expansion. The integration by part analogue to (III.1.23) and the vertical expansion derive some low momentum fields  $\phi_{\rho-1}$  which (together with the true external fields, of scale -1) should be considered as external variables at this stage ( $\phi_{\rho-1}$  is usually called the “background field” at this stage).

The  $\rho-1$  cluster expansion is then performed, with respect to fields  $\phi^{\rho-1}$ , the gaussian measure  $d\mu^{\rho-1}$  and the lattice of cubes of  $\mathbf{D}^{\rho-1}$ . Remark that this may result in propagators connecting cubes of  $\mathbf{D}^{\rho-1}$  already connected (through scale  $\rho$ ). Then we perform the  $\rho-1$   $t$ -decoupling expansion (in which all the remaining former external fields  $\phi_{\rho-1}$  get decomposed into  $\phi^{\rho-1}$  or  $\phi_{\rho-2}$ ). We get a formula which is a sum over  $(\rho-1)$ -polymers which are maximal connected objects made of disjoint sets of cubes of  $\mathbf{D}^{\rho-1}$  and of the  $\rho$ -polymers previously generated. The low momentum fields  $\phi_{\rho-2}$  are the new external set of variables. We iterate the process, generating  $i$ -polymers, until scale 0 is reached.

At the end of the horizontal and vertical cluster expansion the partition function is expressed directly as a sum over vacuum 0-polymers of the type (III.1.15), with hardcore constraints over the whole phase space. The pressure could be computed by a global Mayer expansion of these constraints. For Schwinger functions, we would have only the true external fields (with index -1) left as external variables of the 0-polymers and we can choose the last Mayer expansion adapted to the precise quantities we want to compute (normalized or truncated Schwinger functions). This Mayer expansion generates configurations made out of these connected 0-polymers, joined through “-1 Mayer links”. These Mayer links, which force the polymers to overlap, may also be considered as new connections, so that connected functions are given by a single such configuration with the prescribed set of true external variables according to (III.1.46).

It is time now to define the analogue of almost local subgraphs. Just as before with simple graphs it is useful to consider, for a given 0-polymer and a given slice  $i$ , the set of all  $i$ -subpolymers  $G_k^i$ , which are defined inductively as the maximal  $i$ -polymers (made of cubes of  $\mathbf{D}^i$  and of  $(i+1)$ -subpolymers) which are

connected through connections of indices higher or equal to  $i$ . Again the  $G_k^i$ 's will be connected later together by lower scale connections and they have therefore tree structure.

A convergent 0-polymer is the obvious generalization of the convergent assignments for graphs at the beginning of Section II.3. It is a 0-polymer such that any almost local subpolymer  $G_k^i$  has at least 5 external variables (low momentum fields or true external fields) hooked to it. (As for graphs, in the last slice this condition requires that there at least 5, and by parity in fact 6 true external fields; hence strictly speaking we have no completely convergent 0-polymers unless we consider Schwinger functions with at least 6 external arguments; but this problem may be technically circumvented if necessary by performing a single global subtraction.) It may be interesting to remark again that any  $G_k^i$  which contains a cube  $\Delta \in \mathbf{D}^i$  with  $t_\Delta \neq 0$  (open gate) automatically satisfies the convergence condition because the  $t_\Delta$  expansion has been pushed to 4<sup>th</sup> order so that in the remainder at least 5 low momentum legs have been generated.

The reader should be aware that many technical features of the expansion are optional. In particular the rules we have just described are simple and systematic but they are far from optimal from the point of view of “minimal expansionism”. For instance, having read section III.1, the reader might object that the brute force formulas (III.2.2)-(III.2.12) reminds him of the “pair of cubes” expansion and that presumably they do not give the minimal way to decouple slices in the vertical direction. This guess is correct; just like the “pair of cubes” expansion sometimes builds redundant horizontal loops, the vertical expansion defined above can build redundant vertical connections, i.e. link vertical regions with much more than the minimal number of fields necessary for convergent power counting. This phenomenon can be avoided and there is a more economical way in terms of expansion steps to decouple vertical regions, but as usual it is of a more inductive character. Instead of introducing one parameter  $t_\Delta$  for each cube of the  $\rho$ -th slice, we may introduce a single parameter  $t$  for each of the regions of  $\Lambda$  which have been connected by the ordinary horizontal cluster expansion in the  $\rho$  slice performed before [dCdVMS]. In the previous formalism this is equivalent to equate, for all  $\rho$ -polymers  $P$  the parameters  $t_\Delta$  for  $\Delta \in P$  to a single parameter  $t_P$ , and to perform the Taylor expansion (III.2.12) for all  $t_P$ 's rather than all  $t_\Delta$ 's. In the same vein one should, in the  $\rho - 1$  horizontal cluster expansion, introduce interpolating  $s$  parameters which test the coupling through  $C^{\rho-1}$  of the blocks of cubes connected together through the previous connections, instead of testing blindly between all cubes of  $\mathbf{D}^{\rho-1}$ . At the end of this  $\rho - 1$  cluster expansion, we have a gas of  $\rho - 1$  polymers  $P'$  and new parameters  $t_{P'}$  are defined for the  $\rho - 1$  vertical expansion through the collapse of the corresponding set of  $t_\Delta$  parameters,  $\Delta \in \mathbf{D}^{\rho-1}$ , and so on. This optimized version has many advantages, in particular the superficially convergent  $i$ -polymers (i.e those with 5 external legs  $\phi_{i+1}$  or more) do coincide with the remainder terms in the vertical expansion, so that in a sense they stand out more clearly; also if phase space was used for practical computations it would be presumably optimal to expand in this minimal way. Finally it seems a general rule that expansions with minimal decoupling are always better from the point of view of preserving positivity. In models with more marginal positivity than  $\phi^4$ , such as the Gross-Neveu model in three dimensions of section III.4B, the optimized

expansion becomes truly necessary. Here for simplicity we choose to stick to the systematic rule (III.2.2). The convergence proofs that we provide below also apply with some modifications to the optimized expansion, for which we refer the reader to [dCdVMS].

We propose to study the convergence of the phase space expansion in the situation which corresponds to Weinberg's uniform theorem (II.1.9) or more precisely to its generalization (II.2.1). In this case every contribution has to be convergent from the point of view of power counting. Hence the analogue of Weinberg's uniform theorem in the phase space expansion is:

**Theorem III.2.1 Convergence of the phase space expansion in the case of completely positive power counting**

The sum over all convergent polymers of the phase space expansion is absolutely convergent, uniformly in  $\rho$  and  $\Lambda$ , provided  $g$ , the coupling constant, is small enough; it has a limit as  $\rho \rightarrow \infty$  and  $\Lambda \rightarrow \infty$  which one might call the convergent piece of the  $\phi_4^4$  theory. (Of course this limit is not a field theory and depends strongly on the particular rules we have taken for slicing, etc., ..., because every polymer which violates the convergence condition has been arbitrarily set to 0).

**Proof** The proof does not reduce only to a tedious exercise in patching together everything we know by now about cluster expansions and tree combinatorics in the strange “ $d+1$ ” dimensional phase space. It involves a new difficulty which did not appear at all until now, which is the possible proliferation of low momentum fields. The solution to this problem is called the “domination” of low momentum fields, and it uses in a crucial way the positivity of the interaction. For a better understanding we provide now an informal discussion of the problem and its solution.

**B) Overview of the domination problem**

First one should remember that although the initial functional integral has been quite chopped by the phase space expansion, the subpolymers still contain true functional integrals, hence are not simply ordinary graphs. Of course this was already true for the single scale cluster expansion of the last section. But in that case to bound the functional integral which was of the form

$$\left( \prod \phi_j \right) e^{\int -g\phi^4} d\mu(\phi) \quad (\text{III.2.13})$$

we applied simply a Schwartz inequality to separate  $\prod \phi_j$  from  $e^{-\int g\phi^4}$ ; in other words we used positivity of  $e^{-g\int \phi^4}$  to bound it by 1, and we bounded  $(\prod \phi_j)$  by integrating it with  $d\mu(\phi)$ ; this is what we call “gaussian integration” of the fields  $\phi_j$  produced by the expansions. The corresponding bound is similar to what would be obtained in perturbation theory.

This is all right for a single-scale model. However, when many scales are present, and a given vertex  $v_\Delta$  in  $\Delta \in \mathbf{D}^i$  produced by the horizontal and vertical expansions of scale  $i$  has both high momentum fields  $\phi^i$  and low momentum fields  $\phi^j$  hooked to it,  $j \ll i$ , we can use “gaussian integration” for the  $\phi^i$  fields, but it could be unwise to use it for the  $\phi^j$  fields. Indeed, we would lose the correct power counting or generate a fraction of the factorial divergence of perturbation theory, depending on which point of view is used. We will discuss these two points

of view to understand well the problem; then we will present the solution, again discussing it from both points of view. The domination problem is indeed tricky and worth being discussed thoroughly.

At scale  $i$  the computation of the cluster expansion involves functional derivations which can be expressed as  $C^i(x, y) \frac{\delta}{\delta \phi^i}(x) \frac{\delta}{\delta \phi^i}(y)$ ,  $x \in \Delta$ ,  $y \in \Delta'$ ,  $\Delta, \Delta' \in \mathbf{D}^i$ . These functional derivations  $\frac{\delta}{\delta \phi^i}(x)$  can derive vertices whose three other fields are not of scale  $i$  but of a lower scale, by the obvious analogue of (III.1.23):

$$\frac{\delta}{\delta \phi^i} \cdot e^{-g \int_{\Delta} \left( \sum_{j=1}^{\rho} \phi^j \right)^4} = -4g \int_{\Delta} \left( \sum_{j=1}^{\rho} \phi^j \right)^3 e^{-g \int_{\Delta} \left( \sum_{j=1}^{\rho} \phi^j \right)^4} \quad (\text{III.2.14})$$

Also the vertical expansion in the parameter  $t_{\Delta}$  typically creates such vertices, both in the remainder term  $R^{(4)}$  of (III.2.12) (at least five low momentum fields are derived) and in the term  $I^{(4)}$  of (III.2.12), which is taken at  $t_{\Delta} = 0$ ; remark however that in this last case the number of such fields is at most 4. Let us consider what happens if like in the previous, single scale model, we use a Schwarz inequality to evaluate the corresponding derived low momentum fields of index  $j \ll i$  by gaussian integration.

We have to remark that the gaussian piece available in  $\Delta$  for  $\phi^j$  will be much weaker than for  $\phi^i$ . To see this let us write intuitively the gaussian measure like in (I.3.1) as  $e^{-\int \phi(p^2 + m^2)\phi} D\phi$  in terms of an ill defined Lebesgue measure  $D\phi$ . Then the slicing cutoffs tell us that typically  $p^2$  in the slice  $j$  is of order  $M^{2j}$ , hence much smaller than  $p^2$  in the slice  $i$ , which is typically of order  $M^{2i}$ . Therefore using  $e^{-\int_{\Delta} \phi^j p^2 \phi^j} D\phi^j$  to integrate over  $\phi^j$  is worse than using  $e^{-\int_{\Delta} \phi^i p^2 \phi^i} D\phi^i$ ; by simple scaling it results in a loss of  $M^{(i-j)}$  per low momentum field integrated in this way, in comparison with the estimates of the single scale model (the effect of  $m^2$  is negligible for an ultraviolet problem where  $i$  and  $j$  are big).

The defect of this point of view is that the Lebesgue measure does not exist; nevertheless it leads to the correct conclusion, as we see now by introducing the more rigorous second point of view, in which the gaussian integration over low momentum fields like  $\phi^j$  is rigorously evaluated as a sum of graphs. In this second point of view the domination problem shows up as a piece of the divergence of perturbation theory. We know that gaussian integration leads inescapably to local factorials of the number of fields integrated (see Lemma II.6.2). For fields  $\phi^i$  the corresponding local factorials are harmless because they can be beaten by the volume argument (Lemma III.1.3). Ultimately this is because the cubes of the  $i$ -th cluster expansion are of the correct size for the  $i$ -th momentum slice. But this is no longer true for low momentum fields  $\phi^j$  created by the  $i$ -th cluster expansion; in other words, a lot of  $\phi^j$  fields may be produced in a single cube  $\Delta$  of  $\mathbf{D}^j$ , each of which coming from a different cube  $\Delta'$  of  $\mathbf{D}^i$ , simply because there are  $M^{4(i-j)}$  such different cubes  $\Delta'$  in  $\Delta$ . In Fig.III.2.3, a typical “worst scenario” is pictured: it is a simplified situation, in which each cube  $\Delta' \in \mathbf{D}^i$  sends three low momentum fields (corresponding to a single vertex) in  $\Delta \in \mathbf{D}^j$ .

By our former remarks if  $\Delta \in \mathbf{D}^j$  and  $t_{\Delta''} = 0$  in the  $t$  decoupling expansion for every  $\Delta'' \subset \Delta$ ,  $\Delta'' \in \mathbf{D}^{j+1}$ , only a finite ( $4 \times M^4$ ) number of fields  $\phi^j$  created by higher scales expansions can enter  $\Delta$ . Factorials of such a constant number

are harmless. The worst-case situation of Fig.III.2.3 can occur only if the  $t_{\Delta''}$  parameters for cubes  $\Delta''$  intermediate between the cubes  $\Delta'$  of Fig.III.2.3 and the cube  $\Delta$  are nonzero. In other words the corresponding cubes have to be connected by open gates and belong to the same strongly connected domain. These domains where the interaction is completely undecoupled are therefore truly responsible for the domination problem.

To measure precisely the local factorial effect generated by the gaussian integration of low momentum fields, consider the situation of Fig.III.2.3; let  $n = M^{4(i-j)}$  be the total number of vertices, which is also the number of cubes  $\Delta'$  in  $\Delta$ , since we postulated one vertex per such cube. The total number of fields  $\phi^j$  localized in  $\Delta$  is then  $3 \cdot M^{4(i-j)} = 3n$ . There is no permutational symmetry among the vertices since they lie in different cubes  $\Delta'$ . The gaussian integration of the high momentum  $\phi^i$  fields costs only  $(\text{const})^n$  (by Lemma II.6.2), but the gaussian integration of the  $3n$  low momentum fields  $\phi^j$ , by the same principle, costs  $(\text{const})^n (\frac{3n}{2})!$ . Finally propagators  $C^j$  have a scaling factor  $M^{+2j}$  instead of  $M^{2i}$  for  $C^i$ , and this creates in the estimate a relative bonus of order  $M^{-2(i-j) \cdot \frac{3n}{2}} \simeq (\text{const})^n [(\frac{3n}{4})!]^{-1}$ . Hence the total contribution is of order  $C^n (\frac{3n}{4})!$ , consistent with the intuitive picture that “three-fourths” of the ordinary divergence of perturbation theory has been developed, because three fourths of the  $n$  fields are packed in a single cube of the size corresponding to their frequency. It is also consistent with the heuristic prediction of the first point of view: a loss of  $M^{(i-j)} = n^{1/4}$  per low momentum field means here a total loss of  $n^{(1/4) \cdot 3n}$ , corresponding to the factor  $C^n (\frac{3n}{4})!$ .

A physicist might be shocked by Fig.III.2.3, and remark that “one high momentum leg and three low momentum legs” is a somewhat surprising contribution which does not obey momentum conservation, hence might be suppressed in some way or another. However momentum conservation is partly incompatible with the principle of phase space localization, so that to implement this idea is difficult. A vertex summed over a small cube corresponds in momentum space to a convolution which can bring momenta of order the inverse size of the cube. To exploit momentum conservation at the lower scale it is therefore essential to integrate the corresponding vertices over cubes of a sufficiently large size, hence to free them from the constraints of localization at the higher scale. Moreover technically it requires also cutoffs which preserve momentum rather well. In spite of these difficulties this idea is interesting and will be used in the next section for the Borel summability results. But it is clearly not enough to solve alone the domination problem. We see indeed that even in the case (certainly allowed by momentum conservation) where each vertex would have only one or two low momentum legs, a divergent factor (respectively  $(\frac{n}{4})!$  or  $(\frac{n}{2})!$ ) would still appear.

Therefore we cannot escape the conclusion that in such cases the functional integration over low momentum fields should not be done using solely the decrease of the gaussian measure.

The only other alternative available is to make better use of the decrease of the  $\phi^4$  interaction itself, and immediately everything falls into its proper place. Returning to the first point of view, we saw that  $e^{-\int_{\Delta} \phi^j p^2 \phi^j}$  was much weaker than  $e^{-\int_{\Delta} \phi^i p^2 \phi^i}$ , because in the first case  $p^2 \simeq M^{2j}$  and in the second  $p^2 \simeq M^{2i}$  so that in the first case the typical size of  $\phi^j$  is  $M^{2i-j}$ , much larger than the typical

size  $M^i$  of  $\phi^i$ . But this is no longer true if we use  $e^{-g \int_{\Delta} (\phi^j)^4}$ , which is just as good as  $e^{-g \int_{\Delta} (\phi^i)^4}$ ; both gives a typical size  $g^{-1/4} M^i$  to the field (since the volume  $|\Delta|$  is  $M^{-4i}$ ). In other words the  $\phi^4$  interaction is scale invariant, hence much better in this respect than the gaussian.

In the more rigorous second point of view, where we consider the large number  $(3n)$  of fields  $\phi^j$  which accumulate in a single square  $\Delta$  to give  $(\frac{3n}{2})!$  by Wick theorem or Lemma II.6.2, we know that a quartic integration gives only half as much factorials as a gaussian one; this is nothing but the formula

$$\int x^n e^{-x^2} dx \sim (c_1)^n \left(\frac{n}{2}\right)!, \quad \int x^n e^{-x^4} dx \sim (c_2)^n \left(\frac{n}{4}\right)! \quad (\text{III.2.15})$$

Therefore quartic interaction will result in a  $(\frac{3n}{4})!$  (times  $g^{-\frac{3n}{4}}$  to account for the coupling constant in front of  $\phi^4$ ). Combining this with the relative bonus  $(\frac{3n!}{4})^{-1}$  which is unchanged since it comes from the scaling of  $\phi^j$  (now of order  $M^j$  instead of  $M^i$ ) we obtain  $(\text{const})^n g^{-\frac{3n}{4}}$ , or more precisely  $(\text{const})^n g^{+\frac{n}{4}}$  if the factor  $g^n$  in front of the  $n$  derived vertices is included. This is perfectly summable if  $g$  is small enough.

We realize that it was unreasonable to expect earning a full small factor  $g$  for each derived vertex, when possibly only one fourth of the fields of these derived vertices were of the correct scale! We must abandon this pretension and accept to earn only a more reasonable  $g^{1/4}$  factor, still sufficient of course for convergence purposes. This is done in two steps: gaussian integration of the high momentum fields, a positive operation in which we earn  $g^{1/4}$  per such field, and a bound on the unwanted low momentum fields which compares them back to the exponential from which they were derived; this second operation is neutral from the point of view of estimates (no significant gain or loss except a constant per field).

This sketchy solution of the “domination” problem requires as a crucial point not only the positivity of the interaction (which was required already for the simple existence of the functional integral in a single cube) but also its decay at large  $\phi$  (see (III.2.15)).

We have completed our tour of the problem and its solution, but we want to treat a particular simple example with more care to uncover some subtleties. In particular the reader might have remarked already at this stage that the exponential of the interaction contains interpolating parameters  $t$  which are not necessarily the same than those of the low momentum fields to be dominated, and might worry about that. He might also ask how to separate in practice the high momentum fields from the lower ones in order to apply the domination principle. We address these questions now.

The functional integral within a strongly connected set (set of cubes connected by open gates  $t_{\Delta} \neq 0$ ) may be bounded as

$$|\int ABCP d\mu(\phi)| \leq P \left[ \int B^2 d\mu(\phi) \right]^{1/2} \sup_{\phi} |AC| \quad (\text{III.2.16})$$

where  $A$  is the exponential of the interaction,  $B$  is the product of all high momentum fields (i.e. whose index equal or greater to the index of the expansion step which created them);  $C$  is the product of the low momentum fields,

and  $P$  is the product of the explicit factors produced by the cluster expansions: explicit propagators,  $s$  and  $t$  parameters and coupling constants. The interaction in the exponential is a sum of type (III.2.4), each term having a particular  $t$  dependence and a particular smearing function. We must check that the low momentum fields in  $C$  have  $t$  dependence and smearing functions compatible with the desired bounds. For simplicity let us use sharp characteristic functions. Two main technical points must be understood, namely the smearing of the low momentum fields and the reconstruction of the  $t$  dependence. Both problems are solved in the same way, by introducing some new, harmless set of “high momentum fields”. Let us show in some detail how this is done for the simplest possible example, a model with only two slices,  $i$  (the higher one) and  $j < i$  (the lower one). The field is  $\phi = \phi^i + \phi^j$  and we consider a cube  $\Delta \in \mathbf{D}^i$ . We want to apply a bound of type (III.2.16) to the concrete case:

$$A = e^{-g \int_{\Delta} (\phi^i + t\phi^j)^4 + (1-t^4)(\phi^j)^4} \quad (\text{III.2.17})$$

$$\text{“}BC\text{”} = \int_{\Delta} \phi^i(x)(\phi^j(x))^3 d^4x \quad (\text{III.2.18})$$

The exponential of the interaction is smeared with the characteristic function of  $\Delta$  (in simple words: integrated in  $\Delta$ ). The vertex is also smeared in  $\Delta$ . However we cannot use this smearing function directly to compute the supremum in (III.2.16), precisely because we want to separate the high momentum field from the low momentum field. (This is an elementary point, but one that we stress because it took us some time to understand it...). We have therefore to factorize first (III.2.18), which is not yet a true product of the form  $BC$  (hence the quotes around  $BC$ ). This is done by smearing each low momentum field, hence by writing:

$$\phi^j(x) = \left(\frac{1}{|\Delta|}\right) \int_{\Delta} \phi^j(y) d^4y + \delta\phi^j(x) \quad (\text{III.2.19})$$

The fluctuation field  $\delta\phi^j$  may be written as an integral of derived fields  $\partial\phi^j$ :

$$\delta\phi^j = \left(\frac{1}{|\Delta|}\right) \int_{\Delta} d^4y \int_0^1 dt (y-x)^\mu \partial_\mu \phi^j(x + t(y-x)) \quad (\text{III.2.20})$$

Such a  $\partial_\mu \phi^j$  field has an improved power counting; in this two-slice model it can be considered as a high momentum field and integrated with the gaussian, since the factor  $(y-x)^\mu \partial_\mu$  gives a bonus of  $M^{j-i}$ , exactly like in renormalization (see sect. II.2, after II.2.8). This is because the  $\partial$  acting on  $\phi^j$  gives  $M^j$  and  $|x-y| \leq \sqrt{4M^{-i}}$  since both  $x$  and  $y$  belong to the same cube  $\Delta$  of  $\mathbf{D}^i$ . This factor  $M^{j-i}$  compensates the bad factor  $M^{i-j}$  lost before, so that it is legitimate from the point of view of estimates to consider  $\delta\phi^j$  in this simple two-scale example as a high momentum field. In the general multiscale case, this is no longer true, because we need to gain more than  $M^{j-i}$ , to allow for summation over  $j$ . Hence we have to apply again formula (III.2.19), but now to  $\partial_\mu \phi$ . The remainder contains  $\partial\partial$  fields which can be now really considered as high momentum fields (for them we gain  $M^{2(j-i)}$ ). We also get smoothed  $\partial\phi$  fields which can be dominated if, with some hindsight, we include in  $A$  a small piece  $e^{-\int \partial_\mu \phi \partial^\mu \phi}$  of the gaussian,

treated as an interaction. This rule corresponds to the fact that we are safe only when the propagator becomes summable in  $x$ -space, see Theorem III.1.1b. The corresponding formulas are given in the next section for the infrared case, so we do not elaborate further here on this subtlety.

Developing the third power of (III.2.18) we obtain a sum of terms now factorized, one of which is e.g.  $B'C'$ , with:

$$B' = \int_{\Delta} \phi^i(x)(\delta\phi^j(x))d^4x, \quad C' = \left( \int_{\Delta} \phi^j(y)d^4y \right)^2 \quad (\text{III.2.21})$$

$C'$  is not yet in a form suitable to take a supremum with  $A$ . Nevertheless we see that there cannot be any serious difficulty because the interaction in (III.2.17) interpolates between  $(\phi^j)^4$  and  $(\phi^i + \phi^j)^4$ , and because each low momentum field  $\phi^j$  may be written as  $\phi^i + \phi^j - \phi^i$  at the cost of introducing some new high momentum fields. More precisely we write each  $\phi^j$  as

$$(1-t)\phi^j + (\phi^i + t\phi^j) - \phi^i \quad (\text{III.2.22})$$

and expand. Again the  $\phi^i$  fields are high momentum fields which lead to an other redefinition  $B''$ ,  $C''$  of  $B'$ ,  $C'$ , and in  $C''$  we have two possible types of smeared low momentum fields left,  $(1-t)\phi^j$  and  $(\phi^i + t\phi^j)$ . We use for them a Hölder inequality, taking into account that  $|\Delta| = M^{-4i}$ :

$$(1-t)\left(\frac{1}{|\Delta|}\right) \int_{\Delta} \phi^j(y)d^4y \leq M^i \left[ (1-t^4) \int_{\Delta} (\phi^j(y))^4 d^4y \right]^{1/4} \quad (\text{III.2.23})$$

Similarly:

$$\int_{\Delta} (\phi^i + t\phi^j)(y)d^4y \leq M^i \left[ \int_{\Delta} (\phi^i + t\phi^j)^4(y)d^4y \right]^{1/4} \quad (\text{III.2.24})$$

and we conclude by the bound  $x^{1/4}e^{-x} \leq \text{const}$ . Hence each low momentum field dominated produces (up to a constant) a factor  $g^{-1/4}M^i$ , as expected.

The general rules for domination are derived from this simple example, and shown now in more detail. Several technically different solutions exist, e.g. for the definition of the vertical decoupling expansion, for the smearing of the low momentum fields, the use of smoothed or sharp characteristic functions for the lattice cubes, the shapes of the slicing cutoffs, etc.... These technicalities are important but they may obscure the simple mechanism at work; so we suggest that the reader stops for a while to make sure that he has a good intuitive understanding of this mechanism at least in the two-scale case, before going on.

### C) Proof of Theorem III.2.1

We prove Theorem III.2.1 in a slightly different context, where domination is easier. We write the initial gaussian measure  $d\mu$  as  $e^{-\epsilon \int \partial_\mu \phi \partial^\mu \phi} d\mu'$  with  $\epsilon$  a small constant (we will see that  $\epsilon = g^{1/2}$ , where  $g$  is the bare  $\phi^4$  constant, is a convenient choice). The slicing and horizontal cluster expansion is performed with respect to  $d\mu'$ , and the small quadratic piece is treated as an interaction, using

formulas (III.2.7-9-11)\*. The slicing is such that the sliced propagator  $C^i$  has scaled exponential decrease (II.1.7) or scaled power law decrease

$$|C^i(x, y)| \leq c \cdot M^{2i} (1 + M^i |x - y|)^{-r} \quad (\text{III.2.25})$$

with  $r$  as large as will be necessary. We assume (III.2.25); the case of exponential decrease is easier.

In order to perform the functional integral according to the principle (III.2.16), we should give now the precise rules to determine in the general case what is a low and a high momentum field; this includes, as sketched in the above, the smearing operation and a  $t$ -dependence reconstruction for the low momentum fields in order to rewrite them in a form directly suited for domination. We explain these steps now in full detail.

We define first the raw high and low momentum fields. Then we prepare the raw low momentum fields for domination, i.e. we separate them into true low momentum fields suited for domination plus some new high momentum fields. At the end of this process we have the true or final low and high momentum fields to which formula (III.2.16) is applied.

A field which is not in  $A$ , the exponential of the interaction, must have been derived from it by a cluster expansion step or by a  $t_\Delta$  derivation at some scale  $i$ . In the first case it is hooked to a vertex to which an explicit propagator  $C^i$  of the  $i$ -th cluster expansion is also hooked; in the second case, it is hooked to a vertex in  $\Delta$  which has both fields of indices smaller than  $i - 1$  and bigger than  $i$ . By definition such a field derived at stage  $i$  is either a low momentum field  $\phi_{i-1}$ , i.e. a sum over frequencies  $i - 1, i - 2$  etc... which we do not develop, or it is a high momentum field, which we can in contrast decompose systematically into a sum over some frequencies  $j \geq i$  of sliced fields  $\phi^j$ . We distinguish between the index  $j$  of such a high momentum field and the index  $i \leq j$  (called its production index) of the expansion which produced it.

To be produced at scale  $i$  is the end of the story for a high momentum field, but it is not for a raw low momentum field  $\phi_{i-1}$ , because these ones can be rederived later. When derived by a cluster expansion of scale  $j < i$ , a field  $\phi_{i-1}$  simply disappears into the production of a new explicit propagator  $C^j$  and we need not worry about it any longer. But since the field  $\phi_{i-1}$  becomes really a function  $\phi_{i-1}(t)$  according to (III.2.3), it can be also derived by vertical  $t$  expansions of scale  $j < i$ . Such a derivation selects the piece  $\phi_{j-1}(t)$  (and destroys the higher frequencies, with smaller indices). The last index  $j$  of this type is then called the production index of the low momentum field, and the cube  $\Delta \in \mathbf{D}^j$  to which it belongs is called its production cube. We can essentially forget about the initial  $i$ .

Strongly connected domains are useful to describe in a precise way the range of frequencies of a high or low momentum field produced by some expansion step.

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\* Strictly speaking this changes slightly the definition of what is the convergent piece of the phase space expansion in Theorem III.2.1. This is a minor point since this definition depends also of the arbitrary shape of the slicing cutoffs, and Theorem III.2.1 (which in fact holds in a rather general context) was only intended as a pedagogical introduction to the study of the convergence of phase space expansions.

More precisely a low momentum field with production cube  $\Delta \in \mathbf{D}^j$  is equal either to  $\sum_{k=l(\Delta)}^{j-1} \prod_{k < m \leq j} t_m \phi^k$  or to  $\sum_{k=l(\Delta)}^{j-1} \prod_{k < m < j} t_m \phi^k$ , where  $l(\Delta)$  is the index of the largest cube (coresponding to lowest frequency) in the strongly connected domain to which  $\Delta$  belongs, depending to whether the corresponding vertex was produced by a horizontal or vertical cluster expansion step and whether the low momentum field was rederived later (in which case we have always the first form). Since  $t_j \leq 1$ , we can bound systematically each such field by the first case (putting if necessary the additional  $t_j$  in  $P$ ).

A high momentum field with production cube  $\Delta \in \mathbf{D}^j$  and position  $x$  is equal to a sum

$$\sum_{k=j}^{h(x)} \prod_{k < m \leq h(x)} t_m \phi^k \quad (\text{III.2.26})$$

where  $h(x)$  is the index of the smallest cube (hence with highest scale) containing  $x$  in the strongly connected domain to which  $\Delta$  belongs (see Fig.III.2.2).

Let us prepare the low momentum fields for domination. We can immediately perform the reconstruction of the dependence in the  $t$  parameters, which generalizes (III.2.22). We write systematically such a field with production index  $j$  as:

$$\sum_{k=l(\Delta)}^{j-1} \prod_{k < m < j} t_m \phi^k = \sum_{k=l(\Delta)}^j \prod_{k < m \leq j} t_m \phi^k + (1 - t_j) \sum_{k=l(\Delta)}^{j-1} \prod_{k < m < j} t_m \phi^k - \phi^j \quad (\text{III.2.27})$$

The first two fields are considered as true low momentum fields with production index  $j$ , production cube  $\Delta \in \mathbf{D}_j$  and range  $[j, l]$ . The last field  $\phi^j$  is considered a high momentum field with production index  $j$ .

The fields with derivatives, such as those created by derivations of the interaction term  $e^{-(\epsilon/2) \int \partial^\mu \phi \partial_\mu \phi}$  are treated exactly in the same way, with again the proviso that the derivatives are performed before the  $t$ -dependence (see (III.2.7)).

Finally we have to perform a smearing operation so that low momentum fields can be bounded directly by the interaction  $A$ . The formulas are analogues to (III.2.19-20), but (as announced) we push them one step further. We replace (III.2.19) by:

$$\phi(x) = \left( \frac{1}{|\Delta|} \right) \left[ \int_{\Delta} \phi(y) dy + \int_{\Delta} (x - y)^\mu \partial_\mu \phi(y) dy \right] + \delta\phi(x) \quad (\text{III.2.28})$$

$$\delta\phi(x) \equiv \left( \frac{1}{|\Delta|} \right) \int_{\Delta} dy \int_0^1 dt (1 - t)(x - y)^\mu (x - y)^\nu \partial_\mu \partial_\nu \phi(x + t(y - x)) \quad (\text{III.2.29})$$

Similarly we write:

$$\partial_\mu \phi(x) = \left( \frac{1}{|\Delta|} \right) \int_{\Delta} \partial_\mu \phi(y) dy + \delta \partial_\mu \phi(x) \quad (\text{III.2.30})$$

$$\delta \partial_\mu \phi(x) \equiv \left( \frac{1}{|\Delta|} \right) \int_{\Delta} dy \int_0^1 dt (x - y)^\nu \partial_\mu \partial_\nu \phi(x + t(y - x)) \quad (\text{III.2.31})$$

We apply these formulas to every low momentum field with  $\Delta$  the production cube of the low momentum field. In this way we replace each low momentum

field by one or two corresponding fields smeared in a cube of the production scale  $j$ , plus a fluctuation field which has at least two derivatives acting on it. This fluctuation field is considered a high momentum field of production index  $j$ .

The definition of high and low momentum fields is now completed and we can apply formula (III.2.16). It remains to check that the corresponding bounds are sufficient for convergence. Let us check first the bound for smeared low momentum fields produced in  $\Delta \in \mathbf{D}_j$ . The fields of type  $\phi$  produced by (III.2.28) are dominated using the quartic piece of  $A$ . By (III.2.27) they may be of two types:

$$\left(\frac{1}{|\Delta|}\right) \left[ \int_{\Delta} dy \sum_{k=l(\Delta)}^j \prod_{k < m \leq j} t_m \phi^k \right] \quad (\text{III.2.32})$$

or

$$\left(\frac{1}{|\Delta|}\right) \int_{\Delta} dy (1 - t_j) \sum_{k=l(\Delta)}^{j-1} \prod_{k < m < j} t_m \phi^k \quad (\text{III.2.33})$$

We apply the Hölder inequality (III.2.23-24) to these fields. Now by definition of  $l(\Delta)$   $t_{j-1}, \dots, t_{l(\Delta)}$  are not put to 0 by the vertical expansion. Hence the exponential of the interaction contains the required pieces to bound these fields (this is true no matter whether  $t_j$  itself is 0 or not). Therefore we can conclude, using the bound  $(x^{1/4})^m e^{-x} \leq (m/4)!$ . The overall scale factor associated to such a low momentum field with production index  $j$  is  $g^{-1/4} M^j$  (or  $g^{-1/4} M^{j-1}$ , which is the same up to a constant). In particular the power counting factor,  $M^{-j}$ , is the same (up to a constant) as for a  $\phi^{j-1}$  field in perturbation theory; and the local factorial  $(m/4)!$  is better than the local factorial  $(m/2)!$  that gaussian integration of such fields  $\phi_{j-1}$  would have produced.

Let us discuss now the almost similar case of smeared fields of type  $\partial_{\mu} \phi$  produced by (III.2.28) or (III.2.30). They are dominated using a Schwarz inequality similar to (III.2.23-24) with 4 replaced by 2; then the factor  $e^{-(\epsilon/2) \int \partial^{\mu} \phi \partial_{\mu} \phi}$  in  $A$  contains the pieces required to bound these fields, using  $(\sqrt{x})^m e^{-x} \leq (m/2)!$ . It remains to check that the couplings are all right. Domination in such a way produces a factor  $\epsilon^{-1/2} M^{+2j-2}$  for such a field. But  $\epsilon^{-1/2} = g^{-1/4}$ . When some  $\partial_{\mu} \phi$  legs hooked to an ordinary vertex of the  $\phi^4$  type are dominated in this way, there are at most three of them, and domination consumes  $g^{-3/4}$ ; it remains a small factor  $g^{1/4}$  for the vertex. Similarly when the vertex is of the  $\partial^{\mu} \phi \partial_{\mu} \phi$  type, there is at most one leg dominated in this way, which consumes  $\epsilon^{-1/2}$ , hence again a small factor  $\epsilon^{1/2} = g^{1/4}$  remains for the vertex (this explains our choice of  $\epsilon = g^{1/2}$ ).

We have to check also that the fluctuation fields produced at scale  $j$  give again the same factors than a field  $\phi^{j-1}$  with gaussian integration. This is because such a field of frequency  $k < j$ , apart from inessential factors, has two derivatives and two factors  $|x - y|$  of explicit size  $M^{-2j}$ . When integrated with the gaussian (we may use a Schwarz inequality to separate these fields from the rest), they gives a scale factor  $M^{3k}$  and a local factorial  $(m(\Delta)/2)!$  where  $m(\Delta)$  is the number of such fields per square of  $\mathbf{D}^k$ . The total scale factor for these fields is therefore  $M^{3k-2j}$  hence  $M^{j-1} M^{-3(j-k)}$  (up to a constant). A factor  $L^{-1/2} \equiv M^{-2(j-k)}$  per field is what is needed to transform the local factorials at scale  $k \sqrt{m_{\Delta}}!$  into the product of local factorials at scale  $j$  (or  $j-1$ )  $\prod_{\Delta' \subset \Delta, \Delta' \in \mathbf{D}^j} \sqrt{m_{\Delta'}}!$ , because

$\sum_{\sum m'_\Delta = m_\Delta} m_\Delta! / \prod m_{\Delta'}! = L^{-m_\Delta}$ . The last factor  $M^{-(j-k)}$  per field allows to sum over  $k$  at fixed  $j$ . Hence again the effect on the estimates is the same as if these fields were fields  $\phi^{j-1}$ .

The conclusion is that the net effect of the domination of low momentum fields  $\phi$  or  $\partial\phi$  is therefore bounded (up to a constant per field) by replacing first each such field with production index  $j$  and production cube  $\Delta$  by a field  $\phi^{j-1}$  or  $\partial\phi^{j-1}$  localized in  $\Delta$ , next replacing each coupling constant of each vertex by  $g^{1/4}$  instead of  $g$  or  $\epsilon$  and finally performing gaussian integration over all fields.

We show now that the sum of the corresponding contributions in the convergent case (and for a fixed set of at least six external fields) is a convergent series provided  $g^{1/4}$  is small.

It remains to explain the effect of integrating the well localized high momentum fields with the gaussian measure, according to (III.2.16).

Given a convergent 0-polymer containing the external variables, we have bounded its amplitude by a big sum:

$$\sum_T \sum_G \sum_P \sum_\mu \sum_V \sum_W A_{T,G,P,\mu,V,W}$$

This sum has to be performed over:

- the tree shapes  $T$  which give the inclusion structure of the  $i$ -subpolymers (called  $G_k^i$ ,  $k = 1, \dots, l(i)$  by analogy with part II),
- the subpolymers  $G = G_k^i$ ; each  $G_k^i$  is by definition an  $i$  polymer made of connected cubes of  $\mathbf{D}^i \cup \mathbf{D}^{i+1} \dots \cup \mathbf{D}^\rho$ ; remark that its support at scale  $i$ ,  $G_k^i \cap \mathbf{D}^i$ , called also  $S_k^i$ , may be empty; this is analogous to the possibility in section II that subgraphs  $G_k^i$  could have no line at scale  $i$ ,
- the procedures  $P$  which connect together these polymers; this includes whether  $s$  or  $t$  derivatives hook to already derived fields or create new vertices (hook to the exponential),
- the momentum attributions  $\mu$  of scales to the high momentum fields (recall that these fields are decomposed over their allowed range into slice fields). This sum must be performed before the sum over  $P$ , otherwise we do not know exactly which vertices and fields have been created by the expansion,
- the positions  $V$  of the vertices created by the expansion,
- the Wick contractions  $W$  due to gaussian integration applied to all fields (after the low momentum fields with production index  $j$  have been replaced by  $\phi^{j-1}$  fields, as explained above).

Let us consider a square  $\Delta$  in  $\mathbf{D}^i$ , and a vertex  $v$ .  $\Delta$  is called the localization cube of  $v$  if the position of  $v$  is in  $\Delta$  and the highest field hooked to  $v$  has scale  $i$ . Note that if  $\Delta'$  is the production cube of  $v$  we may have  $\Delta \neq \Delta'$  but  $\Delta$  is necessarily in the strongly connected domain of  $\Delta'$ . We should call  $n(\Delta)$  the number of vertices with localization cube  $\Delta$ . We call  $f(\Delta)$  the number of fields of scale  $i$  whose position is in  $\Delta$ .

Then we perform the Wick contractions. At each scale  $i$ , the integration over the fields  $\phi^i$  with  $d\mu^i(s)$  factorizes into subintegrations associated to each  $G_k^i$  (more precisely each  $S_k^i$ ), hence does not modify the definition of the polymers. Using some fraction of the scaled decay of the propagator as in Lemma II.6.2 or Lemma III.1.2, we transform the corresponding sum over Wick contractions into

local factorials. Finally we integrate over the positions of each vertex. We find that the contribution  $A_{T,G,P,\mu} = \sum_V \sum_W A_{T,G,P,\mu,V,W}$  is bounded by:

- a constant per field (or per vertex, which is the same),
- a factor  $g^{-1/4}$  per vertex,
- a factor  $M^i$  per field  $\phi^i$  or (by (III.2.25)) per half propagator  $C^i$ ,
- a factor  $M^{-4i}$  for each vertex localized in a cube of scale  $i$ , due to the corresponding volume of integration,
- a factor  $\sqrt{f(\Delta)!}$  per cube  $\Delta$  by the local factorial principle,
- a scaled decrease (see (III.2.25)) for each explicit propagator  $C^i$  of the  $i$ -th horizontal cluster expansion.

Let us call  $N(G_k^i)$  the number of external legs of  $G_k^i$  and  $l(G_k^i)$  the set of explicit propagators  $l$  of scale  $i$  in  $G_k^i$ ; such a propagator joins two cubes  $\Delta(l)$  and  $\Delta'(l)$  of  $S_k^i$ . The bound may be summarized as:

$$\begin{aligned} |A_{T,G,P,\mu}| &\leq \prod_{\Delta} \sqrt{f(\Delta)!} (c.g^{-1/4})^{n(\Delta)} \prod_{i,k} M^{-N(G_k^i)} \\ &\quad \prod_{l \in l(G_k^i)} [1 + M^i \cdot \text{dist}(\Delta(l), \Delta'(l))]^{-r} \end{aligned} \tag{III.2.34}$$

Let us show now that we can perform the sum over momentum attributions and compatible procedures and get rid of the local factorials  $\sqrt{f(\Delta)!}$  using a piece of the decrease in  $\prod M^{-N(G_k^i)}$ . For this purpose it is convenient to define  $r(v)$ , the range of a vertex, as its length in the vertical direction, hence the difference between its localization scale and the scale of the lowest field attached to it.

Since our polymers are convergent,  $N(G_k^i) > 5$ , hence  $N(G_k^i) \geq 4 + N(G_k^i)/5$  and we can replace the power counting factor  $\prod_{i,k} M^{-N(G_k^i)}$  in (III.2.34) by  $\prod_{i,k} M^{-4}$  times  $\prod_v M^{-r(v)/5}$ , which gives an exponential decay in the vertical direction for each vertex. Holding the localization scale of each vertex fixed, we can perform the sum over the scale of each leg attached to this vertex, hence over momentum attributions with a third of this vertical decay (at the cost of a constant per vertex, see section II.1). Then let us control the sum over procedures. For each  $s$  or  $t$  derivative we may choose by a factor 2 whether it derives the exponential or hooks to a field already produced. In the second case we have a factor  $f(\Delta)$  to pay to choose this field. In this way we reach a bound of the same type than (III.2.34) but with  $\sqrt{f(\Delta)!}$  replaced by  $f(\Delta!)^{3/2}$  (this is a very crude estimate). It remains to beat such local factorials by the decrease of the cluster propagators.

Let us call  $f^j(\Delta)$  the number of fields  $\phi^i$  localized in  $\Delta \in \mathbf{D}^i$  with production index  $j$ . By our rule for domination we have  $j \leq i+1$ , hence  $f(\Delta) = \sum_{j \leq i+1} f^j(\Delta)$ . We define also  $f_k(\Delta) = \sum_{j \leq k} f^j(\Delta)$  (hence  $f_{i+1}(\Delta) = f(\Delta)$ ). Furthermore, for  $\Delta' \in \mathbf{D}^j$  and  $j \leq i$ :

$$\sum_{\Delta \subset \Delta'} f^j(\Delta) \leq 15 + 3d(\Delta') \tag{III.2.35}$$

where  $d(\Delta)$  is the coordination number of  $\Delta$  in the  $i$ -th cluster expansion (the analogue of  $d_i$  in section III.1). This is because the  $t_\Delta$  expansion produces at most 15 high momentum fields and each horizontal cluster expansion step which

links  $\Delta$  to an other cube creates at most 3 high momentum fields. This is also true for low momentum fields, hence we have also:

$$f^{i+1}(\Delta) \leq \sum_{\Delta' \in \Delta^{i+1}, \Delta' \subset \Delta} (15 + 3d(\Delta')) \quad (\text{III.2.36})$$

Applying repeatedly the binomial argument, then applying (III.2.36) with the remark that there are  $M^4$  cubes of  $\Delta^{i+1}$  in  $\Delta$ , we can write

$$\begin{aligned} f(\Delta)! &\leq \prod_{j=0}^{i+1} 2^{f_j(\Delta)} f^j(\Delta)! \leq [c^{f(\Delta)}] 2^{\sum_{j=0}^i (i-j)f^j(\Delta)} \\ &\quad \prod_{\Delta' \in \Delta^{i+1}, \Delta' \subset \Delta} [d(\Delta')!]^3 \prod_{j=0}^i f^j(\Delta)! \end{aligned} \quad (\text{III.2.37})$$

where  $c$  is some ( $M$ -dependent) large constant.

By (III.2.35) we have, , for  $\Delta' \in \mathbf{D}^j$  and  $j \leq i$ :

$$\prod_{\Delta \subset \Delta'} f^j(\Delta)! \leq c^{\sum_{\Delta \subset \Delta'} f^j(\Delta)} [d(\Delta')!]^3 \quad (\text{III.2.38})$$

hence, combining (III.2.37) and (III.2.38), if  $f \equiv \sum_{\Delta \in \mathbf{D}} f(\Delta)$  is the total number of the fields (and  $c' = c^2$ ) we get:

$$\prod_{\Delta \in \mathbf{D}} f(\Delta)! \leq (c')^f \prod_{\Delta \in \mathbf{D}} [d(\Delta)!]^6 \prod_i \prod_{\Delta \in \mathbf{D}^i} 2^{\sum_{j=0}^i (i-j)f^j(\Delta)} \quad (\text{III.2.39})$$

Provided  $M$  is chosen large enough, the factor  $\prod_i \prod_{\Delta \in \mathbf{D}^i} 2^{\sum_{j=0}^i (i-j)f^j(\Delta)}$  is beaten by an other third of our vertical decay  $\prod_v M^{-r(v)/15}$ ; and the factor  $\prod_{\Delta \in \mathbf{D}} [d(\Delta)!]^6$  is beaten by the horizontal decay of the propagators by lemma III.1.3. (We do not look for optimal bounds; in particular we may avoid to take  $M$  large by a more careful analysis, left to the reader). Finally  $f$ , the total number of fields, is at most four times  $n$ , the total number of vertices, and  $f$  is certainly larger than  $|G|$ , the total number of cubes in the last polymer  $G^0$ . Therefore for any large constant  $K$ , choosing  $g$  small enough we have

$$K^f g^{-n/4} \leq g^{-|G|/8}. \quad (\text{III.2.40})$$

We have still at our disposal a piece of the horizontal decay of the propagators and the last third of our vertical decay. Hence at this stage we have reached the bound on  $A_{T,G} = \sum_{P,\mu} A_{T,G,P,\mu}$ :

$$|A_{T,G}| \leq g^{-|G|/8} \prod_{i,k} M^{-4} \prod_{l \in l(G_k^i)} [1 + M^i \text{dist}(\Delta(l), \Delta'(l))]^{-r/4} \prod_v M^{-r(v)/15} \quad (\text{III.2.41})$$

We have to perform finally the sum over  $T$  and  $G$  of  $A_{T,G}$ . At this stage it is convenient to introduce  $\bar{G}$  which is obtained from  $G$  by “filling the vertical holes”: a cube  $\Delta$  belongs to  $\bar{G}$  if there are two cubes  $\Delta'$  and  $\Delta''$  in  $G$  with  $\Delta' \subset \Delta \subset \Delta''$ . Similarly we define  $\bar{G}^i$  as the set of cubes of  $\bar{G}$  with scales  $\geq i$ , and  $\bar{G}_k^i$  as the

corresponding maximal connected components (see Fig.III.2.4). The advantage is that the set of cubes of  $\bar{G}_k^i \cap \mathbf{D}^i$ , called  $\bar{S}_k^i$ , is never empty, in contrast with  $S_k^i$ . With a factor  $2^{|\bar{G}|}$  we can find  $G$  from  $\bar{G}$ . We may use our last piece of vertical decay  $M^{-r(v)/15}$  to obtain a small factor for each cube in  $\bar{G}$  but not in  $G$ ; this factor will be as small as necessary provided  $M$  is chosen large enough. In particular we may use half of this small factor to perform the sum over  $G$  at fixed  $\bar{G}$  and still retain a small factor per cube of  $\bar{G}$ .

Hence there is a constant  $\epsilon$ , as small as we want provided  $M$  is taken large enough and  $g$  small enough (in this order), such that our bound contains a small constant  $\epsilon$  per cube of  $\bar{G}$ . The final sum may be written  $\sum_{T,\bar{G}} A_{T,\bar{G}}$ . To evaluate the bound that we have at this stage for  $A_{T,\bar{G}}$ , there is still a small difficulty in the way one should describe the horizontal decrease still available. With the power law decrease (III.2.25) the “replica trick” (II.1.22) is no longer possible. Furthermore the cluster expansion at stage  $i$  does not exactly connects together the cubes of  $\bar{S}_k^i$ , but only some subsets  $\bar{S}_{k,m}^i$  such that  $\cup_m \bar{S}_{k,m}^i = \bar{S}_k^i$ . The various  $\bar{S}_{k,m}^i$  are connected together through connections of scale strictly higher than  $i$ . If we call  $T_{k,m}^i$  the tree built by the  $i$ -th cluster expansion between the cubes of  $\bar{S}_{k,m}^i$ , and define the corresponding scaled tree decay:

$$d_{k,m}^i \equiv \prod_{(\Delta, \Delta') \in T_{k,m}^i} (1 + M^i \cdot \text{dist}(\Delta, \Delta'))^{-r/4} \quad (\text{III.2.42})$$

we have the bound (to be compared to the left hand side of (III.1.28)):

$$|A_{T,\bar{G}}| \leq \epsilon^{|\bar{G}|} \prod_{i,k} M^{-4} \sum_{\bar{S}_k^i; Ext \in \bar{S}_1^0} \sum_{\{\bar{S}_{k,m}^i\}; \cup_m \bar{S}_{k,m}^i = \bar{S}_k^i} \sum_{T_{k,m}^i \text{ tree in } \bar{S}_{k,m}^i} \prod_{i,k,m} d_{k,m}^i \quad (\text{III.2.43})$$

where  $Ext \in \bar{S}_1^0$  recalls that the cubes of  $\bar{G}$  at the last scale have to contain the external variables which break the last translation invariance and provide anchoring to the whole construction.

We can picture the connections of the  $\bar{S}_{k,m}^i$  by the graph  $\Gamma$  of Fig.III.2.5. Each value of  $(i, k, m)$  is a node at height  $i$ . It is no longer a tree, but it has to be connected. From  $\Gamma$  we can recover  $T$  by looking at the connected components of  $\Gamma$  above any given height  $i$ . In  $\Gamma$  there is a node at level 0 which contains external variables and which we choose as the root of the construction. Then a partial ordering relation on the nodes of  $\Gamma$  is defined in terms of the minimal number of lines of  $\Gamma$  to reach the root from the node.

We can delete some lines of  $\Gamma$  to form a tree  $T'$  such that for each node the number of steps in  $T'$  to reach the root is this minimal number of steps in  $\Gamma$ . Of course  $T$  can be recovered from  $T'$  rather than from  $\Gamma$ , so we divide our sum over  $T$  into smaller sums indexed by  $T'$ . Then we organize the sum over each  $\bar{S}_{k,m}^i$  in the order given by the tree  $T'$ , in a way which is very similar to the Mayer expansion in section III.1. We start from the nodes which are farthest from the root in  $T'$  and integrate inductively over the corresponding  $\bar{S}_{k,m}^i$ . For each node  $(i, k, m)$ , using the corresponding factor  $d_{k,m}^i$ , we can perform the sum over the trees  $T_{k,m}^i$  and over the position of all the cubes of  $\bar{S}_{k,m}^i$ , save one,  $\Delta_{k,m}^i$  which is kept fixed, and which overlaps with the ancestor of  $(i, k, m)$  in  $T'$ . By the scaled analogue

of Lemma III.1.4, this summation results simply in a constant to the power the number of cubes in  $\bar{S}_{k,m}^i$ .

We have finally to give a rule for the choice of  $\Delta_{k,m}^i$ . At each node we may choose by a factor 2 whether the tree  $T'$  goes up or down in phase space; this factor will be beaten provided  $\epsilon$  in (III.2.43) is small enough. If the ancestor is  $\bar{S}_{k',m'}^{i+1}$  with scale  $i+1$  (i.e. one goes up in the diagram), it is enough to know the cube of  $\bar{S}_{k',m'}^{i+1}$  with which  $\Delta_{k,m}^i$  overlaps, because this cube fixes  $\Delta_{k,m}^i$ . If the ancestor is  $\bar{S}_{k',m'}^{i-1}$  (i.e. one goes down in phase space) one has to choose the cube of  $\bar{S}_{k',m'}^{i-1}$  which contains  $\Delta_{k,m}^i$ , but it is not enough; one has also to pay a factor  $M^4$ , the number of cubes of  $\mathbf{D}^i$  in a cube of  $\mathbf{D}^{i-1}$ . But the corresponding factors  $M^4$  compensate exactly the factors  $M^{-4}$  in (III.2.43), because for each value of  $(i,k)$  there is a *single* value of  $m$  such that one goes down at node  $(i,k,m)$  in  $T'$ .

Once this is realized, the choice of the cube in the ancestor is similar to the problem solved in section III.1 for the single scale Mayer expansion (equations (III.1.39-43)). If we have a coordination number  $d_j$  at node  $j$  in  $T'$ , this leads to a factor  $|\bar{S}_{k,m}^i|^{d_j-1}$  and there is an overcounting symmetry factor  $(n!)^{-1}$  if the summations are made independent over the  $n$  nodes of the tree; by Cayley's theorem,  $1/n!$  is changed into  $\prod 1/(d_j-1)!$ . We can bound  $\sum_{d_j} |\bar{S}_{k,m}^i|^{d_j-1}/(d_j-1)!$  by  $e^{|\bar{S}_{k,m}^i|}$ , which itself is bounded using the small constant per cube of  $\bar{G}$  in (III.2.43). Hence the full sum is bounded by a convergent geometric series if  $M$  is large enough and  $g$  is small enough.

Let us also remark that we could also use the mechanism of convergence of the phase space cluster expansion Theorem III.2.1 as the starting point for a global Mayer expansion in which we could quotient out the “convergent partition function” defined as the sum over vacuum polymers which have no divergent subpolymers at any scale. But it is time now to come to more realistic situations involving renormalization.

### III.3. The effective phase space expansion and infrared $\phi_4^4$

#### A. Model and results

In this section we add renormalization to the phase space expansion (more precisely *useful* renormalization), to obtain finally a tool sufficiently powerful for a non-perturbative investigation of some properties of  $\phi_4^4$ . In particular we give the results on the construction of the critical  $\phi_4^4$  with fixed ultra-violet cutoff, or infrared  $\phi_4^4$ , based on infrared asymptotic freedom [FMRS5], and we include also some corresponding triviality results on ultraviolet  $\phi_4^4$ .

In the case of infrared  $\phi_4^4$  we start with a bare theory of type (I.3.1) in a finite box  $\Lambda$  with a fixed ultraviolet cutoff and a mass counterterm:

$$d\nu(\phi) = Z^{-1} \cdot e^{-\lambda \int_{\Lambda} \phi^4 - (1/2) \int_{\Lambda} m^2 \phi^2} d\mu_C(\phi) \quad (\text{III.3.1})$$

where  $d\mu_C(\phi)$  is a gaussian measure with fixed ultraviolet cutoff. It is convenient to number the scales in this infrared problem in an order opposite to the ultraviolet one, so that the ultraviolet or bare scale has now index 0, and lower and lower momenta have higher and higher indices; it is now the last infrared scale which is called  $\rho$ , as shown in Fig.III.3.1. Remark that there is now a “ceiling” rather than a “floor” in this picture but it would be wrong to consider that the basic picture (Fig.II.1.2) of phase space has been turned upside down; analysis still proceeds from high to low momenta and the definition of almost-local objects is not changed. In summary shifting from an ultraviolet to an infrared problem is like changing the boundary conditions in phase space, but not the structure of the expansion.

The fixed ultra violet cutoff could be of any type, e.g. lattice, exponential or Pauli-Villars type. For a theory with real bare coupling  $\lambda$  any slicing may be used, in particular our favorite exponential slicing (II.1.3) (with the necessary rescaling to adapt it to an infrared problem). But in order to get the Borel summability results we need to construct the theory with  $\lambda$  in the complex disk of Fig.I.5.1, and it is convenient to use a slicing rule with good momentum conserving properties. This eliminates vertices with one high momentum and three low momentum legs whose “domination” is impossible for a small imaginary coupling at the border of the disk (see below). Therefore we choose to use from the beginning\* a cutoff and a slicing rule which in momentum space is smooth but with compact support (this is optimal from the point of view of momentum conservation). We choose a fixed  $C^\infty$  function with compact support  $\eta$  with  $\eta(x) = 1$  if  $|x| \leq 1$  and  $\eta(x) = 0$  if  $|x| \geq 2$ , and use as cutoff:

$$C(x-y) = \frac{1}{(2\pi)^4} \int e^{ip(x-y)} \frac{\eta(p^2)}{p^2} d^4 p \quad (\text{III.3.2})$$

This leads to the natural slicing rule with good momentum conserving properties:

$$C = \sum_{i=0}^{\infty} C^i, \quad C^i = \int e^{ip(x-y)} \frac{\eta^i(p^2)}{p^2} d^4 p \quad (\text{III.3.3})$$

---

\* In [FMRS5] a slicing rule suited for a cluster expansion of the “pair of cubes” type (but not good for momentum conservation) is introduced first, then for Borel summability one turns to an other slicing rule with good momentum conservation; we fear that this may be the source of some confusion.

where  $\eta_0 = \eta$ ,  $\eta^i \equiv \eta_i - \eta_{i-1}$ , and  $\eta_i(t) \equiv \eta(tM^{2i})$ , so that  $C^i$  vanishes if  $|p| \leq M^{-i}$ . This last point is useful for the Borel summability results. However remark that such momentum conserving slicings can be applied to any initial cutoff, at the price of formulas less elegant than (III.3.3); so the method is really general, not limited to cutoffs like (III.3.2). Since  $\eta$  is smooth, the sliced covariance  $C^i$  satisfies the bound (III.2.25) for any large fixed  $r$ .

$m^2$  is the mass counterterm, which is chosen to fix the theory at the critical point, i.e. to lead to a renormalized massless theory; this means that the asymptotic behavior of the two point function will not be exponential decay as in a massive theory, but rather a power-law decay.  $m^2$  is only meant as a generic name for such a parameter, and will turn out to be negative in this problem, so that the associated  $m$  is purely imaginary.

The main result of [FMRS5] is:

**Theorem III.3.1 Existence of infrared  $\phi_4^4$**

For sufficiently small coupling  $\lambda > 0$  the Schwinger functions of infrared  $\phi_4^4$  exist. Their large distance behavior is gaussian up to logarithmic corrections. They are Borel summable functions of  $\lambda$  for  $\lambda \in C_R$  ( $C_R$  being the disk of Fig.I.5.1, and  $R$  a sufficiently large constant).

Similar results (apart from Borel summability) were obtained in [GK2-3].

Some explicit examples of what is meant by gaussian asymptotic behavior were proven in [FMRS5], in particular:

$$S^2(x, y) = \frac{1 + O(\lambda)}{|x - y|^2} [1 + C_1(\lambda, |x - y|)] \quad (\text{III.3.4})$$

$$|C_1(\lambda, |x - y|)| \leq \frac{C_2}{1 + \log(1 + |x - y|)} \quad (\text{III.3.5})$$

$$|S^4(x_1, \dots, x_4)| \leq \int d^4y \prod_{i=1}^4 S^2(x_i, y) \frac{C_3}{1 + \inf_{1 \leq i < j \leq 4} \log(1 + |x_i - x_j|)} \quad (\text{III.3.6})$$

More precise or more general results involving general  $N$ -point Schwinger functions could also be derived from the expansion by explicit computation of the leading terms in the effective phase space expansion.

## B. Overview

The construction that we give in this chapter differs in a substantial way of the one of [FMRS5]; we have tried to incorporate many improvements and hope that it will be simpler to assimilate. We thank Jacques Magnen for collaboration on these improvements. Here is a list of the main differences:

- a) A single cutoff and slicing (III.3.2-3) is used; tree cluster expansions are introduced from the beginning.
- b) Sharp characteristic functions of cubes are used, so that problems of overlap, corridors, etc... no longer occur.
- c) The smearing of low momentum fields and the domination argument is simpler.
- d) We abandon the 3-rd order cluster expansion and modified Mayer expansion which were introduced in [FMRS5] to perform directly full mass renormalization, and return to the simpler (although less elegant) version in which the

counterterm  $m^2$  is computed by a convergent sequence of approximations, rather than by a global formula. This is discussed again below.

- e) The Mayer expansions which allow renormalization cancellations to work are completely explicit.
- f) The proof of Borel summability is made simpler by using momentum conservation at the beginning to eliminate the corresponding troublesome piece of interaction; in [FMRS5] the rules of the expansion itself were modified to restore momentum conservation within the expansion, which is somewhat more complicated.

The improvements  $b$  and  $c$  are possible because we decide to keep a piece of our initial gaussian in the form of an interaction; this will be helpful later for domination purposes. Hence rather than formula (III.3.1), our starting point is:

$$d\nu(\phi) = Z^{-1} \cdot e^{-\lambda \int_{\Lambda} \phi^4 - (\epsilon/2) \int_{\Lambda} \partial_{\mu} \phi \partial^{\mu} \phi - (m^2/2) \int_{\Lambda} \phi^2} d\mu_C(\phi) \quad (\text{III.3.7})$$

Strictly speaking this would not build the theory with bare wave function constant 1, but  $1 + \epsilon$ . To maintain (III.3.4) (with  $O(\lambda)$  instead of  $O(\epsilon)\dots$ ) we have therefore to correct slightly the propagator into:

$$C = \sum C^i, \quad C^i = \int e^{ip(x-y)} \frac{\eta^i(p^2)}{(1-\epsilon)p^2} d^4 p \quad (\text{III.3.8})$$

The constant  $\epsilon$  has to be small, so as to lead to a convergent cluster expansion, but it should not be too small because it will be used for domination. A good choice is  $\epsilon = \lambda^{1/2}$ .

The basic idea of the effective phase space expansion is to mimick the effective perturbation theory of Sect.II.4 by renormalizing in the phase space expansion of the preceding section the almost local polymers with two and four external fields. We give now an informal introduction to this idea, focusing first on the coupling constant renormalization. Renormalization, like in perturbation theory, is performed by subtracting at 0 external momenta; again in phase space this corresponds to hooking all the external fields of a four point polymer to a single point to get the corresponding counterterm. But how is it possible to create such counterterms? This is equivalent to a modification of the value of the bare coupling constant. Is this constant not definitely fixed in (III.3.7)? These are natural questions. The correct answer is that there is of course really no change in the coupling constant, simply it is essential, in order to prove convergence, that we combine together different terms that the expansion produces in order to effectuate crucial cancellations. The rules for these cancellations lead precisely again to the definition of effective constants: hence these effective constants may be considered either as the deep solution to the problem or simply as intermediate tools to exhibit the convergence hidden in the bare formulas.

Let us give an example based on a simplified model with only two slices, two lattices  $\mathbf{D}^0$  and  $\mathbf{D}^1$  and two fields  $\phi^0$  (high momentum) and  $\phi_1$  (low momentum, by our infrared conventions). Remember that almost local subpolymers of  $\mathbf{D}^0$  with two or four low momentum fields in slice 1 must be supported by cubes  $\Delta$  of  $\mathbf{D}^0$  with interpolating parameters  $t_{\Delta}$  set to 0. But after all the interpolating formula (III.2.1) is partly arbitrary; what matters is that at  $t = 1$  the theory

coincides with the initial one (III.3.7) but nothing requires that at  $t = 0$  both decoupled interactions have the same coupling constants. Hence the introduction of the running coupling constant is in fact extremely easy and natural: we replace the interpolation formula (III.2.1) by:

$$\int -\lambda_0(\phi^0 + t\phi_1)^4 - \lambda_1(1 - t^4)(\phi_1)^4 \quad (\text{III.3.9})$$

where  $\lambda_0 \equiv \lambda$  and  $\lambda_1 = \lambda_0 + \delta\lambda$ ,  $\delta\lambda$  being the counterterm, equal to minus the 0 momentum value of the (normalized, connected) four point function corresponding to a theory with only the 0 momentum slice and bare coupling  $\lambda$ . Let us sketch how interpolation (III.3.9) generates now renormalized almost local configurations. Four  $t$  derivations in the term  $I^{(4)}$  in (III.2.7) may now apply to the new term  $t^4\delta\lambda(\phi_2)^4$ ; hence for each cube  $\Delta$  with  $t_\Delta = 0$  one counterterm  $\delta\lambda$  localized in  $\Delta$  is generated. This counterterm should in some sense renormalize the four point polymers produced by the expansion whose “support” contains  $\Delta$ . But how is this possible? There is indeed a difficulty: the counterterm is a universal object (the 0 momentum value of a four point function), and there cannot be any hardcore constraints in its definition. The four point polymers actually created by the expansion have in contrast hardcore constraints with all other polymers produced. To remove hardcore constraints is what the Mayer expansion does. Hence we understand the need, announced in the previous section, to sandwich the cluster and vertical expansions at each scale with a Mayer expansion. This expansion, discussed in subsection D, has to factorize the vacuum polymers (this corresponds to normalization), but it has also to free the two and four point functions from hardcore interactions (except for their *external cubes*, see below).

How does the renormalization transfer convergence to improve the power counting? Of course by a mechanism very similar to the perturbative example discussed in the beginning of Sect.II.2. The main difference is that gradients created by the renormalization subtraction can no longer apply to propagators but must apply to fields: this is in a sense more natural (see the clumsy distinction between the two ends of a propagator in sections II.2-3). For example a four point polymer

$$\int A(x_1, x_2, x_3, x_4)\phi(x_1)\phi(x_2)\phi(x_3)\phi(x_4) \quad (\text{III.3.10})$$

is renormalized by a counterterm

$$\int A(x_1, x_2, x_3, x_4)\phi(x_1)^4 \quad (\text{III.3.11})$$

by writing:

$$\begin{aligned} & \int A(x_1, x_2, x_3, x_4)[\phi(x_1)\phi(x_2)\phi(x_3)\phi(x_4) - \phi(x_1)^4] \\ &= \int A(x_1, x_2, x_3, x_4)\phi(x_1)[\sum_{i=2}^4(\phi(x_1))^{i-2}(\phi(x_i) - \phi(x_1)) \prod_{j < i, j \leq 4} \phi(x_j)] \end{aligned} \quad (\text{III.3.12})$$

The difference of two fields at different arguments, using a Taylor formula, is similar to a derived field  $\partial_\mu\phi$ ; when this object is evaluated by gaussian integration it leads

to the same effect than a gradient on a propagator, hence to the same transfer of convergence discussed at length in sections II.2–3.

In perturbation theory we discovered that mass renormalization is quite different from coupling constant renormalization, in a way being simpler: it does not contain overlapping divergences, hence does not lead to sums over forests, nor does it lead to renormalons effects, provided we work at the level of one-particle irreducible objects. Hence in perturbation theory there is little incentive to pass to the effective expansion as far as mass renormalization is concerned; in fact we used directly full mass renormalization, e. g. in the construction of planar theories (section II.5). Here again perturbation theory is a good guide for what happens in constructive theory, since it turns out that direct full mass renormalization can be performed in the constructive  $\phi_4^4$  problem as well [FMRS5]. We think that it is conceptually interesting to have such a global formula for the mass counterterm. Nevertheless the technical price to pay is heavy, so that in fact this global formula is not very transparent. One has to perform the analysis in phase space at the level of one particle irreducible objects. This requires “3-rd order” cluster expansions and a rather complicated particular Mayer expansion on the chains of one-particle irreducible two point subgraphs, as mentioned in section III.1. In a sense one can argue that the limiting process has been simply hidden in this complicated Mayer expansion.

Therefore here we decided to use a fixed-point argument, less elegant but technically simpler. We construct a sequence of theories and of mass counterterms such that the renormalized mass is closer and closer to 0, and pass to the limit.

In practice we can construct the theory with renormalized mass  $m$  for any  $m \geq 0$ ; this is rather trivial for  $m > 0$  (a theory with both infrared and ultraviolet cutoffs), but what is not trivial is that the case  $m = 0$  (the infrared theory) can be treated so that all convergence estimates will be uniform in  $m$  near  $m = 0$  (hence the approach to the critical point can be studied).

The construction at  $m = 0$  relies in an essential way on the fact that the recursion relation  $\lambda_2 = \lambda_1 + (\delta\lambda)_1$  iterates into  $\lambda_{i+1} = \lambda_i + (\delta\lambda)_i$ , where  $(\delta\lambda)_i = -\beta_i \lambda_i^2 + \text{higher order terms}$ , with  $\beta_i > 0$  and  $\lim_{i \rightarrow \infty} \beta_i = \beta > 0$  (as in section II.5,  $\beta$  is the scale invariant part of the bubble graph). Hence the behavior of  $\lambda_i$  is  $\lambda_i \simeq \frac{\lambda}{1+i\beta\lambda}$ ; the theory is infrared asymptotically free. The wave function constant flow is bounded and small for small  $\lambda$ , like in the planar theory studied in section II.5, because  $\sum \lambda_i^2$  is finite (and  $O(\lambda)$ ).

Finally let us sketch the difficulties associated with the proof of Borel summability. We want our expansion to work also in the case of a complex bare coupling constant  $\lambda$  in the disk  $C_R = \{\lambda, \operatorname{Re}\lambda^{-1} > R^{-1}\}$  of Fig.I.5.1, as should be the case to apply Theorem I.5.1. A new difficulty arises with domination\*, because we can use only the decrease of the real part of the interaction. On the border of the disk  $C_R$  near  $\lambda = 0$  we have  $\operatorname{Re}\lambda \simeq \frac{(\operatorname{Im}\lambda)^2}{R}$  (with in fact  $|\operatorname{Im}\lambda| \leq \sqrt{R \cdot \operatorname{Re}\lambda}$  for  $\lambda \in C_R$ ).

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\* This difficulty is specific of  $\phi_4^4$ ; Borel summability for superrenormalizable theories ([EMS], [GGS], [MS2]) is somewhat simpler, because the coupling constant is dimensioned so that one can trade analyticity in the coupling constant against analyticity in the mass, which is easier to obtain. This is no longer possible for  $\phi_4^4$ . Also in contrast with the superrenormalizable case, it does not seem easy to improve over the Nevanlinna-Sokal domain of analyticity of Fig.I.5.1.

Hence we can dominate at most two low momentum fields hooked to an imaginary vertex and still keep a small factor. We may for instance use the bound:

$$\text{Im}\lambda(\int \phi^4)^{1/2} e^{-\text{Re}\lambda \int \phi^4} \leq c.\sqrt{R} \quad (\text{III.3.13})$$

for some constant  $c$ . But we can no longer dominate three low momentum fields hooked to an imaginary vertex, like in the worst case situation of Fig.III.2.1. Fortunately this type of vertices violates momentum conservation, as remarked already. With our slicing rule (III.3.3) we can even claim that a vertex  $\int_{\mathbb{R}^4} \phi^{i_1} \phi^{i_2} \phi^{i_3} \phi^{i_4}$ , with  $i_1 \leq i_2 \leq i_3 \leq i_4$  is 0 with probability one with respect to the gaussian measure  $d\mu(\phi) = \prod_i d\mu(\phi^i)$  as soon as  $i_2 > i_1 + 1$ . Indeed with probability one the fields  $\phi^i$  have their support in Fourier space out of the region where  $C^i$  is 0, hence they are supported by momenta  $p$  with  $M^{-2(i+1)} \leq p^2 \leq 2M^{-2i}$ , and there is no way that three momenta lower than  $\sqrt{2}M^{-i_1+2}$  add up to a momentum larger than  $M^{-i_1-1}$  (provided  $M \geq 3\sqrt{2}$ , which we assume now). It would seem that we can therefore identify the theory with these vertices to a theory in which they have been suppressed. This is not exactly so because we cannot define functional integrals directly in the infinite volume limit. There is no obvious solution to this subtlety. After comparing various possibilities it seems to us that a reasonably simple solution is to construct the theory in which the imaginary part of the corresponding bare vertices is suppressed, but not the real part (to suppress also the real part does not seem possible, because the corresponding interaction would no longer be positive). The domination problem therefore never occurs, but the corresponding sequence of approximations to the infrared limit are not analytic in  $\lambda$ , so that at the end of the construction one has to give a separate argument to conclude that the limit is analytic in  $\lambda$ . This is however almost obvious: one can check the Cauchy-Riemann equation  $\frac{\partial}{\partial\lambda} f = 0$ . Indeed the Cauchy Riemann operator creates at least one vertex with coefficient  $(\text{Re}\lambda)$  of the non momentum conserving type, and such objects vanish in the thermodynamic limit  $\Lambda \rightarrow \mathbb{R}^4$ . Other alternative solutions would be to keep the unwanted vertices but to use a smoothed characteristic function for our volume cutoff  $\Lambda$ . The unwanted vertices do not completely disappear but they become tiny boundary effects linked to  $\partial\Lambda$ . However one still has to check that they are sufficiently tiny to compensate for the factorials generated by the failure, in their case, of the ordinary domination argument. This is of course true, but painful. Still more complicated perhaps is the solution of [FMRS5]: to modify the expansion rules so as to partly restore translation invariance when such vertices are produced.

In conclusion the bare theory which is our true starting point, when  $\lambda$  is complex is not (III.3.7), but:

$$d\nu(\phi) = Z^{-1} \cdot e^{-\text{Re}\lambda \int_\Lambda \phi^4 - \text{Im}\lambda \int_\Lambda \phi_{MC}^4 - (\epsilon/2) \int_\Lambda \partial_\mu \phi \partial^\mu \phi - (m^2/2) \int_\Lambda \phi^2} d\mu_C(\phi) \quad (\text{III.3.14})$$

where by definition  $\phi_{MC}^4$  (Momentum Conserving) does not contain the momentum violating pieces:

$$\phi_{MC}^4 = \phi^4 - \phi_{MV}^4 ; \quad \phi_{MV}^4 \equiv 4 \sum_i \phi^i \phi_{i+2} \quad (\text{III.3.15})$$

### C) The cluster and vertical expansion

Let us turn now to the details of the expansion. We introduce the orthogonal decomposition of fields  $\phi = \sum \phi^i$  associated to the slicing (III.3.8), and the corresponding lattice of cubes  $\mathbf{D}^i$  of side size  $M^i$ . Recall that now  $\phi_i \equiv \sum_{j=i}^\rho \phi^j$ .

The expansion is performed scale after scale. At scale  $i$  the first part of the expansion is an ordinary cluster expansion among the cubes of  $\mathbf{D}^i$  as described in Section III.1, with respect to the covariance  $C^i$ .

The second piece of the expansion is the vertical decoupling expansion, designed to decouple the fields of frequency  $i$  from the fields with lower frequencies. As in (III.2.7) an operator is applied which computes a fifth order Taylor formula whithin each cube of  $\mathbf{D}^i$  in an interpolating parameter  $t_\Delta$ . Since we have introduced a different interaction for the imaginary and real part of  $\lambda$ , we should give the corresponding interpolation rules in these vertical decoupling parameters  $t_\Delta$ . We introduce the natural generalization of (III.3.15):

$$(\phi_i(t))_{MC}^4 \equiv (\phi_i(t))^4 - (\phi_i(t))_{MV}^4 ; \quad (\phi_i(t))_{MV}^4 \equiv 4 \sum_{j=i}^{\rho-2} [ \prod_{i \leq k < j} t_k ] \phi^j (\phi_{j+2}(t))^3 \quad (\text{III.3.16})$$

where  $\phi_i(t)$  was defined in (III.2.3); we use again the convention that each  $t_k$  stands for  $t_\Delta$ ,  $\Delta$  being the cube of  $\mathbf{D}^k$  to which  $x$  belongs.

The formula which generalizes (III.2.4) and (III.3.9) is:

$$\lambda \phi^4(x) = \sum_{i=0}^{\rho} (1-t_{i-1}^4) [\text{Re} \lambda_i (\phi_i(t))^4 + \text{Im} \lambda_i (\phi_i(t))_{MC}^4 + \text{Im} (\delta \lambda_{i-1}) (\phi_i(t))_{MV}^4] |_{t=1}(x) \quad (\text{III.3.17})$$

with  $\delta \lambda_{i-1} = \lambda_i - \lambda_{i-1}$ . This may seem complicated, but remark that for real  $\lambda$  it reduces to formula (III.2.4). Remark also that we tolerate momentum violating interactions for the counterterms  $\delta \lambda_{i-1}$  because these ones will be dominated easily (they are of order  $\lambda^2$  rather than  $\lambda$ ) and because otherwise some configurations in the expansion would not be properly renormalized; to show that they are in fact harmless (because they violate momentum conservation) would require a separate argument.

When all parameters  $t$  are set to 1, (III.3.16) coincides with the quartic interaction in (III.3.14). The effective parameters  $\lambda_i = \lambda_{i-1} + \delta \lambda_{i-1}$  are defined inductively;  $\delta \lambda_k$  will be defined below as the 0 momentum value of the almost local 4 point configurations at scale  $k$ ; by definition  $\lambda_0 \equiv \lambda$ .

Finally we have to give a rule for the interpolation of the quadratic pieces considered as interaction in (III.3.14) namely  $m^2 \int \phi^2$  and  $\epsilon \int \partial_\mu \phi \partial^\mu \phi$ . We may or may not perform the wave function renormalization (which turns out to be finite). In order to have always exponentially convergent sums, we choose to perform it. Hence we use the interpolations generalizing (III.2.8-9):

$$\sum_{i=0}^{\rho} (1-t_{i-1}^2) (m^2)_i (\phi_i(t))^2(x) \quad (\text{III.3.18})$$

$$\sum_{i=0}^{\rho} (1-t_{i-1}^2) (\epsilon + a_i) [\partial_\mu \phi_i(t) \partial^\mu \phi_i(t)](x) \quad (\text{III.3.19})$$

with the convention  $a_0 = 0$ , and  $(m^2)_i = (m^2)_{i-1} + (\delta m^2)_{i-1}$ ,  $a_i = a_{i-1} + \delta a_{i-1}$ .

It remains to give, at each scale, the precise definition of the counterterms  $\delta m^2$ ,  $\delta a$  and  $\delta \lambda$ .

#### D) The Mayer expansion and the definition of counterterms

At each scale, after each corresponding cluster and vertical expansion has been performed, we have to define and perform a Mayer expansion. As previously stressed, the main conceptual difficulty is that cluster expansions generate polymers (sets of cubes) but Mayer expansion generate configurations, i.e. sequences of polymers joined by Mayer links. Iterating this process leads to sequences of sequences of sequences... and it is not easy to picture the result. This seems unfortunately an intrinsic difficulty, which can be avoided only by hiding it into an induction: we prefer to underline it right from the beginning.

The Mayer expansion that we want to apply is designed to free the vacuum, two and four point functions from hardcore constraints. However there is an obstacle to apply fully this program: since there can be an arbitrarily large number of polymers with four external low momentum fields, if we were to suppress all hardcore constraints between them, there would be the possibility that an arbitrarily large number of low momentum fields accumulate at the same place, and by the local factorial principle, this would lead later to divergences when these fields are estimated.

We should remark also the vertical expansion can produce only *one* counterterm per cube  $\Delta$  with  $t_\Delta = 0$ . Therefore it would be nice if the two and four point functions had all their external legs in a single external cube, and if such external cubes remained disjoint (with hardcore constraints). After playing sometime with this idea we are lead to the conclusion that the Mayer expansion should apply to the vacuum, two and four point functions, but not to their external cubes, i.e. to the cubes which contain their external fields, and that to prepare the function for renormalization: we should separate it into a local and a renormalized piece, in a way such that the local piece has only one external cube, hence will match exactly with the counterterm.

Let us give more precisely the corresponding rules, in the case of the first slice, with index 0. The cluster and  $t$  expansions have produced a set of disjoint polymers with various number of external low momentum fields  $\phi_1$ . The vacuum polymers are called  $V_1, \dots, V_n$ . The two and four point polymers are then separated into a local part and a renormalized part according to formulas analogue to (III.3.12) (we may treat in the same way the polymers with one or three legs but they vanish by parity considerations so let us neglect them). It is better to use fully symmetric formulas, such as (we temporarily use simply  $\phi$  as the notation for the external fields which are really  $\phi_1$ ):

$$\phi(x)\phi(y) = (1/2)[\phi(x)^2 + \phi(y)^2 + [(x-y)^\mu \partial_\mu \phi(x)]^2 + [(y-x)^\mu \partial_\mu \phi(y)]^2] + [\phi(x)\phi(y)]_{\text{ren}} \quad (\text{III.3.20})$$

$$\phi(x_1)\phi(x_2)\phi(x_3)\phi(x_4) = (1/4)\left[\sum_{i=1}^4 \phi(x_i)^4\right] + [\phi(x_1)\phi(x_2)\phi(x_3)\phi(x_4)]_{\text{ren}} \quad (\text{III.3.21})$$

where  $[\phi(x)\phi(y)]_{\text{ren}}$  is an expression symmetric in  $x$  and  $y$  with at least three derivatives acting on  $\phi$ , and  $[\phi(x_1)\phi(x_2)\phi(x_3)\phi(x_4)]_{\text{ren}}$  is an expression with at least one

such derivative. For instance we may take:

$$[\phi(x_1)\phi(x_2)\phi(x_3)\phi(x_4)]_{\text{ren}} = - \sum_{i < j} \left[ \int_0^1 dt (x_j - x_i)^\mu \partial_\mu \phi(x_i + t(x_j - x_i)) \right]^2 \prod_{1 \leq k \leq 4, k \neq i, k \neq j} \phi(x_k) \quad (\text{III.3.22})$$

Using this decomposition for all polymers, we obtain what we call local and renormalized two and four point polymers. The local two and four point polymers are labelled as  $W_1, \dots, W_m$  and have all their external legs at the same point, hence in a well defined cube of  $\mathbf{D}^0$ , called the external cube,  $E_j$  of  $W_j$ . All other cubes in the support of  $W_j$  form what we call the internal domain  $I_j$  of  $W_j$ . The renormalized polymers contain derived fields at interpolated points such that  $x_i = t(x_j - x_i)$  and since the rules for connectivity do not imply that polymers are convex, we must consider in this case that the cube containing the interpolated point is connected to the polymer; this is a new type of link (with decrease obviously similar to the regular decrease of an ordinary propagator) to add to the list of the previous section.

Polymers with more than 4 external legs or renormalized polymers with two and four external legs will be treated similarly from now on because they have good power counting. They are labelled as  $Y_1, \dots, Y_p$ .

We define the Mayer expansion as removing all hardcore constraints involving one vacuum polymer  $V$  or one internal domain  $I_j$ . Hence at the end, hardcore constraints between the  $Y$ 's and the external variables  $E_j$  do remain. More precisely we have a big sum:

$$\sum_{\{V_i\} \{W_j = E_j \cup I_j, E_j \cap I_j = \emptyset\} \{Y_k\}} \prod_{i,j,k} A(V_i) A(W_j) A(Y_k) P_1 P_2 \quad (\text{III.3.23})$$

where the hardcore interactions are divided into  $P_1$  and  $P_2$ :

$$P_1 \equiv \prod_{i \neq i'} e^{-V(V_i, V_{i'})} \prod_{i,j} e^{-V(V_i, W_j)} \prod_{i,k} e^{-V(V_i, Y_k)} \\ \prod_{j \neq j'} e^{-V(E_j, I_{j'})} \prod_{j \neq j'} e^{-V(I_j, I_{j'})} \prod_{j,k} e^{-V(I_j, Y_k)} \quad (\text{III.3.24})$$

$$P_2 \equiv \prod_{j \neq j'} e^{-V(E_j, E_{j'})} \prod_{j,k} e^{-V(E_j, Y_k)} \prod_{k,k'} e^{-V(Y_k, Y_{k'})} \quad (\text{III.3.25})$$

where  $V(X, Y)$ , as in section III.1 is the hardcore interaction, with value 0 if  $X$  and  $Y$  are disjoint and  $+\infty$  otherwise. The Mayer expansion consists in expanding  $P_1$  according to the algebraic formula (III.1.33), but in the inductive way exemplified by (III.1.37-38).  $P_2$  remains unexpanded. The result is expressed as a sum over configurations, i.e. sequences of polymers. The normalization (sum over vacuum configurations) is fully factorized as in (III.1.44). However the computation does not correspond exactly to truncated functions, but rather to a partial truncation at the level of two and four point functions, which furthermore is restricted by the

fact that the truncations relative to the “external cubes” are not performed. Of course the convergence theorems of section III.1 still apply to such situations. The important point is that with this restricted rule, as in the convergent case, the number of external fields  $\phi_1$  in any cube of  $\mathbf{D}^\rho$  cannot get large except when many corresponding propagators hook to distant other cubes, so that the corresponding local factorials remain controlled by the volume effect (lemma III.1.3) \*.

The result of the expansion is expressed as a sum over 0-configurations, i.e. sequences of 0-polymers joined by Mayer links.

Actually there is a detail which is not yet correct, but that according to D. Knuth’s pedagogy we did not address until now. Each time a Mayer configuration is generated by the overlap of two local two point functions both corresponding to mass terms (i.e. both with external fields  $\phi^2(x)$ ), a four point function is generated which is neither local nor renormalized. So we have to apply again to these objects the decomposition (III.3.24). In the new local four point functions defined in this way there is again an external cube and an internal domain, which contains the previous external cube of one of the two point functions whose overlap formed the four point function. Therefore the global definition of external cube and internal domain at the beginning is not exactly correct.

One correct rule (not unique) is to decompose first the two point functions into local and renormalized parts and define the corresponding internal domain and external cubes of the two point functions; then expand all hardcore constraints involving the internal domain of local two point functions. Only when this is finished the four point functions are defined, decomposed into local and renormalized parts (with the corresponding external cube and internal domain for the local parts). Then the remaining constraints, involving vacuum polymers or internal domains, are expanded.

Performed in this way the Mayer expansion factorizes all vacuum polymers and generates configurations with two or four external fields which are either of the local type (with all external fields hooked at the same point) or of the convergent or renormalized type (hence have favorable power counting). Moreover the two and four point local configurations, apart from their single well defined external cube, are completely free of any hardcore constraints with the other objects of the expansion.

Then the cluster and vertical expansions are performed at scale 1. They generate maximal connected objects made of a set of 0-subconfigurations and of a set of cubes of  $\mathbf{D}^1$ . Of course the Mayer links of the 0-th Mayer expansion are taken into account for connectedness. Let us call 1-polymers these maximal connected objects, with some slight abuse of language. We can define the local and renormalized parts of the 2 and 4 point 1-polymers in the same way than before. The external cubes of the local two and four point functions do have hardcore constraints since by definition they belong to  $\mathbf{D}^1$ . We perform then the scale 1 Mayer expansion by expanding all hardcore constraints of type  $P_1$ , and obtain a

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\* It would be possible to remove all the hardcore constraints involving a *two-point* polymer, because when  $n$  such objects accumulate in a given cube, we recover a  $1/n!$  symmetry factor which compensates for the  $n!$  obtained by gaussian integration of  $2n$  fields. This possibility does not extend to higher point functions, so we will not actually use it.

sum over 1-configurations, in which hardcore constraints of the type  $P_2$  remain (between external cubes or 1-polymers of the local type or 1-polymers with more than four external legs).

Iterating this process, for each scale  $i$  we generate  $i$ -polymers and  $i$ -configurations. It remains to show that with the correct definition of the counterterms  $\delta m^2$ ,  $\delta a$  and  $\delta \lambda$ , every  $i$ -subconfiguration in the expansion is renormalized, i.e. every local two and four point subconfiguration disappears exactly (this corresponds to the perturbative rule for the effective expansion that each almost local subgraph is renormalized).

We give the precise definition of the counterterms and perform renormalization in three steps. The  $i$ -th cluster, vertical and Mayer expansions are performed with  $(m^2)_{i+1} = (m^2)_i$ ,  $a_{i+1} = a_i$  and  $\lambda_{i+1} = \lambda_i$ ; they are called the first level expansions (for slice  $i$ ). Then we define first the mass and wave function counterterms (respectively  $(\delta m^2)_i$  and  $\delta a_i$ ) as minus the sum of all local two-point  $i$ -configurations with fixed external cube  $\Delta$  generated at this first level (respectively with external fields  $\phi^2(x)$  and with external fields  $\partial_\mu \phi \partial^\mu \phi(x)$ ). To transform  $[(x-y)^\mu \partial_\mu \phi(x)]^2$  into  $(1/8)|y-x|^2 \partial_\mu \phi \partial^\mu \phi(x)$  is similar to (II.3.44-46) and we do not repeat it here.

In this way we obtain correctly normalized values for these counterterms, which clearly do not contain hardcore constraints with other configurations and are independent of  $\Delta$  (up to boundary effects \*).

The full expansion for slice  $i$  is then recomputed with the new values  $(m^2)_{i+1} = (m^2)_i + (\delta m^2)_i$  and  $a_{i+1} = a_i + \delta a_i$  fed into formulas (III.3.18) and (III.3.19), but still with  $\lambda_{i+1} = \lambda_i$ . The result is called the second level expansion, in which every 2 point local  $i$ -configuration is *exactly* cancelled.

Indeed the case in which a local two point configuration with external cube  $\Delta$  is generated by the former first level expansion cancels exactly against the new case in which two  $t_\Delta$  derivatives have generated exactly one counterterm of the type  $(\delta m^2)_i$  or  $\delta a_i$ , and there is no other possibility for local two point  $i$ -configurations.

To check that the cancellation is exact, one has to remark that the counterterm,  $\delta m^2 \int \phi_{i+1}^2(x) dx$  is completely independent of the fields of slices  $i$  and above. Therefore after the second level cluster and  $t$  expansions but before the second level Mayer expansion the value of a two point  $i$ -polymer with one such counterterm produced in  $\Delta$  reduces to this counterterm times the value of the vacuum polymer to which  $\Delta$  belongs. After the second level Mayer expansion, the vacuum graphs are factorized, and the value of the  $i$ -configuration becomes therefore 1 times the counterterm, hence cancels exactly with the sum of the former (first level) local  $i$ -configurations, as announced. This statement is somewhat tautological. Indeed with the cluster and Mayer expansions of section III.1, we may compute normalized Schwinger functions even in the case where the external fields decouple. When there is a single such external source  $\delta m^2$  in the cube  $\Delta$ ,

\* Because of the boundary effects due to our finite box  $\Lambda$ ,  $\delta m^2$ ,  $\delta a$  and  $\delta \lambda$  have still weak dependence on  $\Delta$ . We decide to neglect this dependence since it has no effect on the results. A way to avoid it completely is to pass to the thermodynamic limit  $\Lambda \rightarrow \mathbb{R}^4$  in each slice successively rather than at the end of the expansion. We leave to the reader to develop this nice possibility.

the unnormalized Schwinger function is obviously  $\delta m^2 \cdot Z$ , where  $Z$  is the normalization. Computed by formula (III.1.44), it is also equal to  $A^T(M) \cdot Z$ , in which  $M$  is the Mayer configuration containing the source. Therefore  $A^T(M) = \delta m^2$ , which is the desired result.

Finally the counterterm  $\delta\lambda_i$  is computed in the same way as minus the sum of all local four point  $i$ -configurations generated at the second level. The new value  $\lambda_{i+1} = \lambda_i + \delta\lambda_i$  is fed into (III.3.17), and the expansion is recomputed for the third and last time. This leads to a third level expansion in which all two and four point  $i$ -configurations are of the renormalized type. Remark that  $\delta\lambda_i$  depends on  $(\delta m^2)_i$ , since two local two point functions (at least one of which is a counterterm) may still combine at second level to give new four point functions (which are automatically decomposed into local and renormalized pieces, hence taken into account for the value of  $\delta\lambda_i$ ). For instance the vertical  $t_\Delta$  expansion, which is a Taylor expansion to fourth order, may create two mass counterterms in a single cube  $\Delta$  with  $t_\Delta = 0$  which is a new kind of four point function. It is for this reason that we introduced these “levels”; in the following we forget about them, and the expansion we speak about is always the final one (third level).

In the case of complex  $\lambda$  there is one additional small technical detail linked to the special form of (III.3.17) which has to be mentioned here: if we study what means  $\phi_{MC}^4$  in formulas (III.3.16-17) we realize that not every four point  $i$ -configuration generated by the expansion has necessarily four external fields of the type  $\phi_{i+1}$ , hence is suitably renormalized by a counterterm of the form  $\phi_{i+1}^4$ . There is indeed the possibility of a (single!) border vertex of a 4 point configuration having only one field in slice  $i$  and three fields below, in which case the three fields become of type  $(\phi_{i+1})^3$  only if we add the missing momentum violating piece  $4\text{Im}\lambda(\phi^i(\phi_{i+2})^3)$ . If we write:

$$4\phi^i(\phi_{i+2})^3 = 4\phi^i[(\phi_{i+1})^3 - 3(\phi_{i+1})^2\phi^{i+1} + 3\phi_{i+1}(\phi^{i+1})^2 - (\phi^{i+1})^3] \quad (\text{III.3.26})$$

we obtain two classes of four point  $i$ -configurations, the normal ones with 4 external fields  $\phi_{i+1}$  and exceptional ones with some of their external fields being  $\phi^{i+1}$ . The obvious rule is then not to take into account these exceptional objects in the computation of  $\delta\lambda_i$ . They should not be considered as true four point  $i$ -configurations. As a consequence they remain unrenormalized at scale  $i$ ; this has no consequence at all, because there is no logarithmic divergence associated to such objects; recall that the need for renormalization of four point almost local objects comes solely from the lack of vertical decrease between their internal and first external scale; but for these exceptional objects, the first external scale is  $i+1$ , next to the last internal scale  $i$ , so there is obviously no divergence associated to them. We will not return again on this small technical detail in the following.

## E) The functional bounds

When the expansion is completed, we write it as a sum of terms

$$\sum \int ABCP d\mu(\phi) \quad (\text{III.3.27})$$

where as in (III.2.16)  $A$  is the exponential of the interaction (including the quadratic pieces),  $B$  is a product of high momentum or well localized fields, and  $C$  is the

product of the low momentum or badly localized fields. As before,  $P$  is the product of all explicit factors (not depending on fields) created by the expansion: this includes the former factors such as propagators derived by the cluster expansions, coupling constants derived by the cluster and vertical expansions and the new  $e^{-V(Y,Y')} - 1$  factors (Mayer links) derived by the Mayer expansions. The principle of the functional bounds is then the same as (III.2.11), namely high momentum fields are integrated with the help of the gaussian measure, and low momentum fields are “dominated”, i.e. bounded with the help of the interaction.

The precise rules to determine what is a low and a high momentum field are the same than in the previous section III.2.C. Beware that in this infrared problem indices run in the opposite way, so that many formulas of the previous section cannot be used literally but summations or inequalities involving indices usually have to be reversed. There is also a minor modification, concerning the  $\text{Im}\lambda$  vertices, in which  $\phi_{i+1}$  must be replaced by  $\phi_{i+2} + \phi^{i+1}$  and  $\phi^{i+1}$  must be treated as a high momentum field. We might even do this systematically for all vertices, so that by convention the low momentum fields would be always of the form  $\phi_{i+2}$ ; the corresponding small technical changes would not be relevant for the main issue of convergence of the expansion.

The  $t$  dependence reconstruction and the decomposition of low momentum fields into smeared pieces plus fluctuation fields is exactly similar to the previous formulas (III.2.26-30). Domination of the smeared low momentum fields proceeds as before, using the  $\phi^4$  interaction for fields and the  $\partial_\mu\phi\partial^\mu\phi$  interaction for derived fields  $\partial\phi$ .

One delicate point however has to be discussed: the coupling constant which equips the vertex to which such a low momentum field  $\phi_{j+1}$  is hooked is  $\lambda_i$  where  $i$  is the *initial* production scale. We will soon check that by asymptotic freedom,  $\lambda_j \simeq \lambda/(1+\lambda\beta.j)$  may be much smaller than  $\lambda_i \simeq \lambda/(1+\lambda\beta.i)$  (see (III.3.28)). The total factor associated to a vertex (apart from power counting factors) is therefore at worst  $\lambda_i\lambda_j^{-3/4} \leq c.\lambda^{1/4}(j-i)$  for some constant  $c$  (the high momentum fields, integrated with the gaussian, do not modify this estimate). We have therefore to find a factor to compensate for the potentially harmful difference  $j-i$  (when this difference gets large). Fortunately power counting in a usefully renormalized expansion is always favorable, so there is exponential decay in the vertical direction; moreover each pair of lines hooked to a vertex has a separate exponential decrease (which is in fact at least  $M^{-|i-j|/18}$ , see (II.1.19)). We can pick a fraction of this decay  $M^{-\epsilon(j-i)}$  and use it to bound the corresponding factors  $(j-i)$ , up to a constant. Remark that different factors  $j-i$  correspond to different pair of fields hooked to a vertex, hence there is no risk of factorials of the renormalon type generated in this way.

Similarly when the vertex is of the  $\partial^\mu\phi\partial_\mu\phi$  type, there is at most one leg dominated. This consumes  $(\epsilon+a_j)^{-1/2}$ , so that at worst  $\sqrt{\frac{\epsilon-|a_i|}{\epsilon+|a_j|}}$  remains for each such vertex. But in the next section we find that  $a_i = \sum_{j=0}^i \delta a_j$  is finite and of order  $\lambda$  for small  $\lambda$ , independently of  $i$  (see (III.3.32)). By our choice of  $\epsilon = \sqrt{\lambda}$ , again a small factor  $\epsilon^{1/2} = \lambda^{1/4}$  (up to a constant) remains for such a vertex.

It remains to discuss the case of mass vertices  $\delta m^2\phi^2$ . It will be shown below (in (III.3.29)) that the correct choice of counterterms is such that  $|(\delta m^2)_i| \leq$

$c \cdot \lambda_i M^{-2i}$ , so that such a vertex, produced in a cube of side  $M^i$  and volume  $M^{4i}$  is neutral just as a  $\phi^4$  vertex; in this case at most one leg is a low momentum leg; its domination with the  $\phi^4$  interaction costs at most  $\lambda_i^{-1/4}$ , hence again a small factor (here at least  $\lambda_i^{3/4}$ ) remains for convergence purposes.

The conclusion is that in every case a small factor per vertex do remain for the convergence of the expansion.

Finally let us examine what happens in the case of a complex coupling  $\lambda$ . As sketched above, domination is all right for all the vertices with at most two badly localized legs, provided the radius of the disk  $C_R$  is chosen small enough. (This includes in particular the case of the quadratic vertices  $\delta m^2 \phi^2$ ). The only dangerous case is the “momentum violating piece”  $\phi_{MV}^4$ , with imaginary coefficient. We choosed formula (III.3.17) precisely so as to eliminate the term  $\text{Im } \lambda \phi_{MV}^4$ . The formula (III.3.17) does contain imaginary vertices of the momentum violating type, but they are counterterms. We will show below that  $\delta \lambda_i \simeq \beta \lambda_i^2$ . Therefore, since  $\lambda^2 \text{Re } \lambda^{-3/4}$  is small for  $\lambda \in C_R$ , these imaginary pieces do not create any problem when bounded by the corresponding  $\text{Re } \lambda \phi^4$  interaction.

The rest of the argument is essentially similar to section III.2.C, and we do not repeat it. Renormalized two and four point functions have indeed the same power counting than a six point function. Let us remark however that after domination and evaluation of the gaussian integrations over high momentum fields, summation is no longer over subpolymers  $G_k^i$  but over subconfigurations, so that the full results of section III.1 (not only the results on the cluster expansion but also theorem III.1.3 on convergence of the Mayer expansion) have to be used in order to control such summations.

A final comment is in order about the external variables  $\phi(x_1), \dots, \phi(x_N)$ . Their treatment depends in some respects of the kind of results we are after. For a simple proof of existence of the correlation functions in the infrared limit, these external fields are smeared against some test functions of the scale 0 (or taken at fixed lattice sites if we work on a lattice). Hence, like in the ultraviolet problem it is convenient to consider by convention these external fields as having index  $-1$ , independently of their true index in the decomposition of each external field  $\phi$  as a sum of  $\phi^i$  fields. But this rule has the opposite effect than in the ultraviolet problem as far as renormalization is concerned: the polymers with e.g. four external fields (and no low momentum fields) have not to be renormalized in the effective expansion. Indeed their power counting is already favorable, because translation invariance is broken at the highest possible scale by the presence of the external fields, which are smeared over well-defined cubes of the smallest size in the problem.

## F) The behavior of the effective constants

The behavior of the effective couplings is studied in the manner of section II.5, i.e. inductively. We may do this in two steps. For the purpose of domination, hence of convergence of the partially renormalized phase space expansion we may prove first inductively a crude bound, then we can study in more detailed the recursion relation which give the flow of effective constants, and prove asymptotic formulas like (III.3.4-6). Some additional work would lead also to the rigorous proof of the more detailed results about logarithmic deviations from mean field

theory in this model which have been derived by the use of the renormalization group. Here we will limit ourselves to check some rather crude bounds, in the form of:

**Lemma III.3.1** There is some constant  $\beta$  and some large constant  $c$  such that:

$$\operatorname{Re}(\lambda_i^{-1}) \geq \operatorname{Re}(\lambda^{-1}) + \beta \cdot i - O(\log i) \quad (\text{III.3.28})$$

$$|\delta(m^2)_i| \leq c \cdot \lambda_i M^{-2i} \quad (\text{III.3.29})$$

$$|\delta a_i| \leq c \cdot (\lambda_i)^2 \quad (\text{III.3.30})$$

These three bounds are derived by looking at the first non trivial term in the definitions of the counterterms. For the coupling constant, this is our friend the bubble graph, which means that the recursion relation is  $\lambda_i = \lambda_{i-1} - \beta \lambda_{i-1}^2 + o(\lambda_{i-1}^2)$ , and this leads to the recursion relation  $\lambda_i^{-1} = \lambda_{i-1}^{-1} + \beta + o(\lambda_{i-1})$ . Starting from some place in  $C_R$  we are driven to the real axis and to the origin, hence (III.3.28) holds.

Similarly the leading graph for the mass counterterm is the single vertex “tadpole”, which is proportional to  $\lambda_i$ ; for the wave function constant the leading graph is  $B_0$  in Fig.II.5.1 and this leads to the quadratic bound (III.3.30). Combining (III.3.29) and (III.3.30) with (III.3.28) we obtain some bounds for the total flow of the mass and wave function terms:

$$\sum_{i=0}^{\infty} |\delta(m^2)_i| \leq O(\lambda) \quad (\text{III.3.31})$$

$$\sum_{i=0}^{\infty} |\delta a_i| \leq O(\lambda) \quad (\text{III.3.32})$$

The reader may wonder why the leading orders are the natural radiative corrections in  $\lambda$  since in this problem there are interactions  $\epsilon \partial_\mu \phi \partial^\mu \phi$  and  $m^2 \phi^2$ . In particular the coefficient  $\epsilon$  is  $\sqrt{\lambda}$  so subgraphs with such vertices seem to dominate. However this is only apparent. We know that if we push the expansion up to a given order  $n$  in  $\lambda$  (hence e.g. up to  $2n$  in  $\epsilon$ ), the graphs with  $\epsilon$  vertices will recombine exactly to change the factor  $1 - \epsilon$  of  $p^2$  in the denominator of each propagator into 1 up to  $o(\lambda^n)$ . Similar remarks apply to the mass vertices, which are shown below to be  $O(\lambda)$ . Hence the renormalization group flows are indeed given by the ordinary radiative corrections.

Remark that the flow of effective constants is simpler in this infrared problem than in the ultraviolet problem of section II.5, because at least for the coupling constant and wave function constant, there are really no renormalization conditions; we start with given values  $\lambda$  and 1, and let them evolve as they want (in fact they go to the gaussian free field theory with some wave function constant close to 1). In particular there is no need of introducing the intermediate values  $\lambda_i^\rho$  as was the case for the  $g_i^\rho$  of section II.5. The only non-trivial renormalization condition is that the renormalized mass should be set to 0. We address this issue now.

## G) The inductive choice of $m_{ren} = 0$

The expansion converges, but what does it construct so far? Not exactly the massless theory, unless the initial mass term  $m^2$  in (III.3.1) is fine-tuned to the correct value. To obtain this, we have to compute

$$m_{ren}^2 = m^2 + \sum_{i=0}^{\infty} (\delta m^2)_i \quad (\text{III.3.33})$$

and to prove that there exists  $m^2$  such that  $m_{ren}^2 = 0$ . Of course the counterterms  $(\delta m^2)_i$  are themselves functions of  $m^2$ , as every other quantity in the phase space expansion. But from (III.3.29) and (III.3.31) we know that the series in (III.3.33) converge, and that  $f_\lambda(m^2) = -\sum_{i=0}^{\infty} (\delta m^2)_i(m^2, \lambda)$  is  $O(\lambda)$ . Therefore for  $\epsilon$  small and  $\lambda$  small enough  $f$  maps the interval  $[-\epsilon, \epsilon]$  into itself. To show that  $f$  has a fixed point, it remains to prove that it is contracting, hence to compute  $\frac{d}{dm^2} f_\lambda$ . The  $\frac{d}{dm^2}$  operator creates a mass term  $m^2 \int \phi^2$  from the exponent of the initial functional integral (III.3.1), (III.3.7) or (III.3.14); it remains to evaluate the effect of such a term on the counterterms  $(\delta m^2)_i$ . This is done with the phase space expansion. The result is simply bounded by  $m^2 \cdot O(\lambda)$ , hence for  $\epsilon$  and  $\lambda$  small the mapping  $f_\lambda$  is a contraction and has a single fixed point  $m^2$ ; this fixed point is  $O(\lambda)$  and at the fixed point we have  $m_{ren} = 0$ , as desired. To derive the asymptotic decay (III.3.4-5) of the propagator at the fixed point is then a simple exercise.

There is one issue which we postpone until now. We remark that the tadpole graph is negative (for positive  $\lambda$ ), hence the leading counterterm  $(\delta m^2)_i$  is positive. Therefore  $m^2$  itself is negative\*. The careful reader may worry about the functional integration with such a term. It is true that one can no longer bound the exponential of the interaction by 1 simply like in (III.1.25). But  $m^2 = O(\lambda)$ , and

$$m_j^2 \equiv m^2 + \sum_{j=0}^{i-1} (\delta m^2)_j = M^{-2i} O(\lambda_i) \quad (\text{III.3.34})$$

since  $\sum_{j>i} (\delta m^2)_j = M^{-2i} O(\lambda_i)$  by an analogue of the analysis above. Hence when we have to bound a functional integral in a cube  $\Delta$  of  $\mathbf{D}^i$  we can bound the dangerous positive mass term in the exponential by half of the negative  $\phi^4$  term, uniformly in  $\phi$ . A prototype of the corresponding inequalities is

$$e^{-\lambda_i \int_{\Delta} \phi_i^4 + (m^2)_i \int_{\Delta} \phi_i^2} \leq K \cdot e^{-(\lambda_i/2) \int_{\Delta} \phi_i^4} \quad (\text{III.3.35})$$

for some large constant  $K$ , since by a Schwarz inequality

$$\int_{\Delta} \phi^2 \leq M^{2i} (\int_{\Delta} \phi^4)^{1/2} . \quad (\text{III.3.36})$$

The large constant  $K$  in the empty cube disappears in the normalization and in non empty cubes is compensated by taking  $\lambda$  still smaller. Half of the large field decrease of the  $\phi^4$  interaction remains available for domination tasks; again some

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\* We recall that  $m^2$  is simply a notation for a parameter which was not assumed to be the square of a real number. We may consider  $m$  to be in fact purely imaginary in this problem.

big constants in the estimates get still bigger, but of course convergence itself is not in danger for  $\lambda$  sufficiently small.

## H) Analyticity and Taylor remainders

To complete the proof of Borel summability, it remains to prove that the correlation functions built by the expansion become analytic in  $\lambda$  when  $\rho, \Lambda \rightarrow \infty$  (this is not true in finite volume, because (III.3.14) is not analytic). Then we have to compute the Taylor remainder at  $n$ -th order in the bare coupling constant  $\lambda$  of these functions and verify a factorial bound (I.5.2).

To verify analyticity, we apply the  $\frac{\partial}{\partial \lambda}$  operator to such a correlation function, as defined by the effective phase space expansion. This operator must act on the only non analytic piece, hence it creates at least one momentum-violating vertex  $\phi_{MV}^4$  swimming somewhere in the phase space expansion. It is therefore similar to a new kind of external source, except that this one is smeared over the whole of  $\Lambda$ , not over a cube of slice 0. We treat this vertex in the same way than other external sources, namely the corresponding fields have scale  $-1$  by convention, hence the divergent polymers which contain them are not renormalized. The limit  $\rho, \Lambda \rightarrow \infty$  is performed as before, and this limit vanishes because such a vertex vanishes when  $\Lambda \rightarrow \infty$ . Of course a true proof requires to check that this vertex does not destroy the convergence of the expansion. If we perform the thermodynamic limit with a smooth  $C_0^\infty$  volume cutoff  $\Lambda(x)$  and link the limit  $\rho \rightarrow \infty$  and  $\Lambda \rightarrow \infty$  by the natural rule  $\int \Lambda(x) dx = M^{4\rho}$ , we find that a momentum violating vertex, integrated in such a volume  $\Lambda$ , gives a very small factor  $M^{-k\rho}$ , where  $k$ , which depends of the exact shape of the smooth function  $\Lambda$ , can be adjusted to any large fixed integer we want. In this way we can compensate for the large volume factor  $M^{4\rho}$  associated to the fact that this new external source is summed over the whole of  $\Lambda$ , instead of being smeared over a small cube of size 0. This proves that in the limit  $\rho \rightarrow \infty$  analyticity is recovered. If we insist on  $\Lambda$  being not smooth but being the characteristic function of some cube of  $\mathbf{D}^\rho$ , we have to take boundary effects into account and the proof is slightly more complicated, but of course the result remains true.

Finally the Taylor remainder bounds are obtained by developing explicitly  $n$  vertices from the exponential before applying the whole phase space expansion. The result converges in the same way than before, but with a prefactor  $n!$  which is simply the large order bounds on perturbation theory detailed at length in part II. Therefore we can apply the Nevanlinna-Sokal theorem and complete the proof of Borel summability of the correlation functions in  $\lambda$ .

## I) Weak triviality

The construction of infrared  $\phi_4^4$  is interesting because it allows to put on a mathematically rigorous basis the famous analysis of this model by the renormalization group [KW]. The logarithmic asymptotic freedom of the model leads to the possibility of a non trivial fixed point of the renormalization group below four dimensions; in its most naive form, this is because the coupling constant is dimensioned out of  $d = 4$ . The corresponding evolution equation for the effective constant  $\lambda_i$  is therefore  $\lambda_{i+1} \simeq \epsilon \lambda_i - \beta \lambda_i^2$ , with  $\epsilon = d - 4$ , and it has a non trivial fixed point at  $\lambda \simeq \frac{\epsilon}{\beta}$ . This fixed point is small for small  $\epsilon$ . These observations are at the basis of the analysis of non-trivial infrared fixed points in powers of  $\epsilon$ .

(the Fisher-Wilson expansion). This expansion leads to reasonably good numerical computations of critical exponents in  $d = 3$  ( $\epsilon = 1$ ).

But the construction of infrared  $\phi_4^4$  sheds also some light on the ultraviolet problem for  $\phi_4^4$ . Theorem III.3.1 can be immediately rephrased as a triviality theorem for ultraviolet  $\phi_4^4$  at weak bare coupling. More precisely:

**Theorem III.3.2: triviality at weak bare coupling**

Let us consider the  $\phi_4^4$  model with fixed (non zero) renormalized mass and wave function constants (for instance computed in the BPHZ scheme), a cutoff  $\rho$  and a bare coupling  $\lambda_\rho$ . There exists some  $\epsilon > 0$  such that the corresponding ultraviolet limit is a free field provided the bare coupling  $\lambda_\rho$  is smaller than  $\epsilon$  for all  $\rho$ .

Of course the weak coupling condition is very restrictive, but it cannot be lifted if we want the phase space expansion to converge. There are more general triviality statements [Aiz][Fr] which apply for any bare coupling but they are limited to a lattice regularization and they operate fully only for  $d > 4$ , with some problems remaining at  $d=4$ . They rely on rigorous inequalities, the only alternative tools we know so far when convergent expansions fail.

On this issue we have to conclude rather sadly that for the moment neither the construction of ultraviolet  $\phi_4^4$  nor a fully general triviality theorem does seem accessible with phase space expansions. This problem remains of more than purely mathematical interest since  $\phi_4^4$  is expected to govern the dynamics of the Higgs particle, and it is not clear how the triviality issue changes when coupled to non-abelian gauge fields and exactly how it affects the corresponding physics.

As a final, more optimistic note, let us mention that the relationship between the bare and the renormalized couplings is not yet fully understood and that some interesting or surprising results may await us in this domain. In particular it should be interesting to study the relationship between the infrared and ultraviolet  $\beta$  functions of  $\phi_4^4$  as is attempted in [Kop]. The infrared  $\beta$  function is defined as the limit, when  $x \rightarrow \infty$  of the derivative of the bare coupling with respect to  $x = \log \frac{\kappa}{m_{ren}}$ , where  $\kappa$  is the ultraviolet cutoff scale and  $m_{ren}$  the renormalized mass. This limit has to be computed by holding the renormalized coupling fixed, and has to be expressed as a power series in the bare coupling. In our notations  $x$  is roughly  $\rho \log M$ . This infrared beta function is the natural infrared analogue of the ultraviolet  $\beta$  function, which is computed by holding the bare coupling fixed and looking at the variation with  $x$  of the renormalized coupling, reexpressed in terms of the renormalized coupling (see the end of section II.5). In the phase space expansion, recall that up to terms in  $(\log M)^p$ ,  $p > 1$ , which reflect the discrete nature of our flow equations, the ultraviolet  $\beta$  function is what we obtain when we compute the discrete change  $\delta\lambda$  in  $\lambda$ . This is certainly a well defined quantity, but phase space expansion expresses it as a power series in all the previous running couplings, and the ultraviolet power series for the  $\beta$  function is obtained only by expressing all these intermediate couplings in terms of the last one, a potentially dangerous operation since it generates useless counterterms usually associated with renormalizations. In contrast the infrared beta function is related to the precise behavior of the renormalized coupling, beyond the obvious  $x$  and  $\log x$  terms, expressed in terms of the bare coupling. In [Kop] it is shown that there is a renormalization scheme in which the ultraviolet  $\beta$  function as a power series in the renormalized coupling is equal to the infrared one as a power series in the bare

coupling, and it is argued that the infrared  $\beta$  function, being related to the infrared Schwinger functions can also be defined by Borel summation. It would be then possible to define the ultraviolet Callan-Symanzik equation non perturbatively at least in this particular scheme. This scheme ressembles the minimal subtraction scheme in dimensional regularization. This gives some weight to the belief that in the dimensional minimal subtraction scheme there are also no renormalons in the  $\beta$  function, a belief which would partly justify its use in the computation of the  $\epsilon$  expansion [BDZ]. Even in the BPHZ scheme the presence or absence of renormalons for the  $\beta$  function is unclear [Kop]. If the  $\beta$  function became well defined a rigorous study of renormalons would be easier, as explained in section II.6. Also it would be quite intriguing that the equations of transformations of  $\phi_4^4$  under a change of scale may be defined non perturbatively. Well defined differential equations do have local solutions; if the corresponding flows are not complete, it is because the solutions explode to infinity in finite time (as we see for the bare coupling). To cure this defect one has usually to reformulate the equations in an other (more compact) space. It is perhaps in this way that one might reach in the future some positive results concerning  $\phi_4^4$ . Let us conclude simply as a philosophical remark that in mathematics non-existence theorems rarely remain the last word on a subject: often a problem with no solution is simply badly formulated and has to wait until the proper formalism in which it does have a solution is found.

### III.4 The Gross-Neveu model

#### A) Two dimensions

The Gross-Neveu model in two dimensions [MW] [GrNe] is a model of fermions with a color index  $N > 1$  and a quartic interaction. The model was introduced in order to study ultraviolet asymptotic freedom, dynamical symmetry breaking and the non perturbative generation of mass in a context simpler than non-abelian gauge theories. The case of a single fermionic field ( $N = 1$ ) is known under the name of the massive Thirring model. This model can be mapped exactly on the bosonic sine-Gordon model and is the subject of an extensive literature, but it is not asymptotically free in the usual sense (its  $\beta$  function vanishes exactly), hence we do not discuss it here.

The massive two dimensional Gross-Neveu model has formal euclidean action:

$$S(\bar{\psi}, \psi) = \bar{\psi}(ia\partial + m)\psi - \lambda(\bar{\psi} \cdot \psi)^2 \quad (\text{III.4.1})$$

where  $\bar{\psi} \cdot \psi \equiv \sum_{a,\alpha} \bar{\psi}_a^\alpha \psi_a^\alpha$ . The letters  $a, b, \dots$  are used for color indices, hence take values from 1 to  $N$  and the letters  $\alpha, \beta, \dots$  are used for spinor indices (which in two dimensions take two values). Pairs of a color and a spinor index such as  $(a, \alpha)$ ,  $(b, \beta)$  are noted  $A, B, \dots$ . Except in the explicit computations of the leading graphs which drive the discrete evolution of our effective constants, the spinor and color indices may be forgotten. The model has a perturbation expansion similar to the one of  $\phi^4$  (more precisely to the one of  $N$ -vector  $\phi^4$ , for which vertices should show the circulation of indices). But it is just renormalizable in two dimensions, just as  $\phi^4$  is in four, because a fermionic propagator decreases at large momenta like  $p^{-1}$ , not  $p^{-2}$ .

We introduce the following conventions for two dimensional  $\gamma$  matrices:

$$\gamma_0 = \begin{pmatrix} i & 0 \\ 0 & -i \end{pmatrix} \quad \gamma_1 = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} \quad (\text{III.4.2})$$

and the usual relations  $\partial = \gamma_0 \partial_0 + \gamma_1 \partial_1$  etc...

The  $2p$  point functions are formally defined as:

$$S_{2p}^{A_1 \dots A_p B_1 \dots B_p}(y_1, \dots, y_p, z_1, \dots, z_p) = \frac{1}{Z} \int \psi_{a_1}^{\alpha_1}(y_1) \dots \psi_{a_p}^{\alpha_p}(y_p) \bar{\psi}_{b_1}^{\beta_1}(z_1) \dots \bar{\psi}_{b_p}^{\beta_p}(z_p) e^{\int \lambda(\bar{\psi}\psi)^2 - \bar{\psi}(ia\partial + m)\psi} \prod_x d\bar{\psi}(x) d\psi(x) \quad (\text{III.4.3})$$

with the usual rules of fermionic (Berezin) integration. We do not know how to define fermionic functional integrals rigorously except by performing them explicitly, which is possible here. Indeed the perturbative expansion of an unnormalized Schwinger function  $S_{2p}^u$  is:

$$S_{2p}^{u, A_1 \dots A_p B_1 \dots B_p}(y_1, \dots, y_p, z_1, \dots, z_p) = \sum_{n=0}^{\infty} \sum_{C_i, D_i} \frac{\lambda^n}{n!} \int_{\Lambda} d^2 x_1 \dots d^2 x_n \left( \frac{y_1 \dots y_p x_1 x_1 x_2 x_2 \dots x_n x_n}{z_1 \dots z_p x_1 x_1 x_2 x_2 \dots x_n x_n} \right)^{A_1 \dots A_p C_1 D_1 \dots C_n D_n} \quad (\text{III.4.4})$$

where the upper variables correspond to the  $\psi$  fields, the lower variables correspond to the  $\bar{\psi}$  fields, and we use the notation:

$$\begin{pmatrix} u_1 \dots u_n \\ v_1 \dots v_n \end{pmatrix}_{B_1 \dots B_n}^{A_1 \dots A_n} \equiv \det(C_{B_i}^{A_i}(u_i, v_j)) \quad (\text{III.4.5})$$

The propagator  $C_B^A$  is diagonal in color space (vanishes unless  $a_i = b_i$ ) and for  $a_i = b_i$  it is equal to  $(-\not{p} + m)/(p^2 + m^2)$ . Let us replace it by a propagator  $C_\rho$  with the same index dependence and our favorite ultraviolet cutoff (of course other ones may be accommodated):

$$C_\rho(p) = C(p)e^{-M^{-2\rho}(p^2+m^2)} = \sum_{i=0}^\rho C^i(p) \quad (\text{III.4.6})$$

$$\begin{aligned} C^i(p) &= C(p)(e^{-M^{-2i}(p^2+m^2)} - e^{-M^{-2(i-1)}(p^2+m^2)}) \quad \text{if } i \geq 1 \\ C^0(p) &= C(p)(e^{-(p^2+m^2)} - 1) \end{aligned} \quad (\text{III.4.7})$$

The indices now run again in the regular way for an ultraviolet problem, and the lattice of cubes  $\mathbf{D}^i$  is again the lattice of cubes with side size  $M^{-i}$ . The sliced propagators are very similar to the bosonic ones; in particular they satisfy for any given fixed  $r$  the bound:

$$|\partial^m C^i(x, y)| \leq K \cdot M^{i \cdot (m+1)} e^{-c \cdot M^i |x-y|} \quad \forall m \leq r \quad (\text{III.4.8})$$

where  $K$  and  $c$  are constants (depending only on  $r$ ), and  $\partial^m$  is any partial derivation of order  $m$ .

With such an ultraviolet cutoff and a finite volume  $\Lambda$ , the series (III.4.4) have in fact an infinite radius of convergence, so that they can be taken as a well defined starting point for the bare theory. Indeed, as we know now, in the cutoff  $\phi^4$  theory the divergence of perturbation theory is a local phenomenon due to the fact that vertices can accumulate in a small spatial region where they become undistinguishable for the propagator. But fermions cannot behave in this way thank to the Pauli principle.

Mathematically this corresponds to the fact that if we develop the determinant (III.4.4) we recover the usual  $\phi^4$  graphs, hence  $n!$  contributions at order  $n$ . However we know that in a determinant there are many changes of signs, so that typically tremendous cancellations can occur. Hence let us start our discussion with a simple bound [IM2], which proves the convergence of (III.4.4) and is also useful for the phase space analysis of the model.

### Lemma III.4.1

For any positive integer  $r$  there exists a constant  $K(r)$  such that, if  $C^i(x_j, y_k)$  is the  $n$  by  $n$  matrix with  $(j, k)$  entries  $C^i(x_j, y_k)$  and for any  $\Delta \in \mathbf{D}^i$ ,  $n_\Delta$  (respectively  $\bar{n}_\Delta$ ) is the number of  $x_j$  variables in  $\Delta$  (respectively the number of  $y_k$  variables in  $\Delta$ ), we have:

$$|\det C^i(x_j, x_k)| \leq [K(r)M^i]^n \cdot \prod_{\Delta} \frac{1}{(n_\Delta!)^r (\bar{n}_\Delta!)^r} \quad (\text{III.4.9})$$

**Proof** Let us give the proof for any dimension  $d$  of space time. The total number of color and spinor pairs is  $2N$ . Let us define  $p = (4r + 2)d$  and divide each cube  $\Delta \in \mathbf{D}^i$  into  $n_\Delta/4Nd^p$  cubes, each of side size  $M^{-i} \cdot (4Nd^p/n_\Delta)^{1/d}$ . Such a new cube is noted  $\Delta_\alpha$ , its center is noted  $z_\alpha$ , and the number of variables  $x_j$  (respectively  $y_k$ ) that it contains is noted  $n_\alpha$  (respectively  $\bar{n}_\alpha$ ). If  $z_\alpha$  is the center of the cube  $\Delta_\alpha$  containing  $x_j$  we apply a Taylor expansion to the propagator  $C(x_j, y_k)$ , writing:

$$\begin{aligned} C(x_j, y_k) &= \sum_{m=0}^{p-1} \frac{1}{m!} \prod_{l=1}^m (x_i - z_\alpha)^{\mu_l} \partial_{\mu_l} C(z_\alpha, y_k) \\ &+ \frac{1}{(p-1)!} \int_0^1 dt (1-t)^{p-1} \prod_{l=1}^p (x_i - z_\alpha)^{\mu_l} \prod_{l=1}^p \partial_{\mu_l} C(x_i + (1-t)(z_\alpha - x_i), y_k) \end{aligned} \quad (\text{III.4.10})$$

In this way each row of the initial determinant is a sum of at most  $d^{p+1}$  vectors. Expanding the determinant, we get a sum of at most  $d^{n \cdot (p+1)}$  determinants. In any of these non-zero new determinants at least  $n_\alpha - 2Nd^p$  of its  $n_\alpha$  rows with arguments  $x_i$  localized in  $\Delta_\alpha$  must be remainder terms in (III.4.10). For each row containing a term with a  $q$ -th order derivative localized in  $\Delta$ , we have using the bound (III.4.8) a net gain in the estimate of the propagator of  $(4Nd^p/n_\Delta)^{q/d}$ . If we expand any such determinant by brute force we obtain therefore the bound

$$|\det C^i(x_j, y_k)| \leq K_1^n \prod_{\Delta} \prod_{\Delta_\alpha \subset \Delta} (1/n_\Delta)^{(n_\alpha - 2Nd^p)p/d} \cdot B \quad (\text{III.4.11})$$

where  $B$  is evaluated as a sum of Feynman graphs without any cancellations taken into account, hence by the local factorial lemma (Lemma II.6.2):

$$B \leq K_2^n M^{ni} \prod_{\Delta} \sqrt{(n_\Delta!)(\bar{n}_\Delta!)} \quad (\text{III.4.12})$$

Since

$$\sum_{\Delta_\alpha \subset \Delta} (n_\alpha - 2Nd^p) = n_\Delta/2 \quad (\text{III.4.13})$$

we obtain

$$|\det C^i(x_j, y_k)| \leq K_1^n K_2^n \prod_{\Delta} (1/n_\Delta)^{n_\Delta(p/2d)-1/2} \sqrt{(\bar{n}_\Delta!)} \quad (\text{III.4.14})$$

and repeating the argument with the columns and taking the geometric mean of the two bounds achieves the proof of the lemma, since  $p = (4r + 2)d$ .

Furthermore we can apply the lemma to  $C_\rho$  instead of  $C^i$ ; if  $\rho$  is fixed to a constant,  $C_\rho$  satisfies the bound (III.4.8) (e.g. with  $i = 0$ ). Therefore the announced result that the bare perturbation expansion for the Schwinger functions with both ultraviolet and finite volume cutoff  $\Lambda$  has an infinite radius of convergence follows. We can integrate the  $x$  and  $y$  positions over the whole of  $\Lambda$ , which give a factor at most  $q^{2n}$  in the bound, if  $q = |\Lambda|$  is the number of unit cubes in  $\Lambda$ . The product over these cubes of e.g.  $n_\Delta! \bar{n}_\Delta!$  is at least  $(2n)!/(2q)^{2n}$ , hence applying the lemma with respect to the cubes of  $\mathbf{D}^0$  with  $r = 1$  we have a series bounded by a series in

$\lambda^n$  of an exponential type, with infinite radius of convergence. Of course the corresponding bounds are very crude, and would not allow a study of the ultraviolet limit. The phase space analysis that we introduce now will precisely allow this by improving vastly over this bound.

Starting with the well defined bare theory (III.4.4-5) and  $C = C_\rho$  we call the bare quantities for  $\lambda$ ,  $m$  and  $a$  respectively  $\lambda_\rho$ ,  $m_\rho$  and  $a_\rho$  as usual.

We apply a phase space expansion similar to the previous case, the determinant in equation (III.4.4) playing the rôle of the exponential of the interaction in  $\phi^4$ ; it is also an expression which cannot be developed completely without loosing the structure of cancellations responsible for convergence. Again the expansion applied to it must preserve most of this structure but develop it sufficiently to show a minimal set of connections in each slice, in order to perform correct power counting analysis and the necessary useful renormalizations (which here again correspond to coupling constant, mass and wave function renormalization).

If we compare further to  $\phi^4$  we realize that the model here is in many respects simpler. The local bound (III.4.9) replaces the domination in an advantageous way because Lemma III.4.1 means that any power of local factorials can be beaten by the Pauli principle, which is therefore in a sense much more powerful than the decrease of  $e^{-\int \phi^4}$  at large  $\phi$ . In fact it is a good intuitive picture to consider Berezinian anticommuting variables as *bounded* variables. Also there is no analogue of the positivity of the  $\phi^4$  interaction, hence no need to preserve it in our interpolation schemes.

The initial constructions of the Gross-Neveu model [FMRS4] [GK4] provided the first renormalizable field theory. It was shown to obey the Osterwalder-Schrader's axioms (for fermions), and the correlation functions were also shown to be the Borel sum of their renormalized perturbation expansion [FMRS4]. Following results on this model include the definition of general irreducible kernels [IM2], a study of the Bethe-Salpeter equation and of two-particle asymptotic completeness [IM4] and the analysis of large momentum properties and the Wilson-Zimmermann short distance expansion [IM3].

In this section we sketch the proof of the main result on the existence of the ultraviolet limit of the model, summarized in:

### Theorem III.4.1 Existence of Gross-Neveu<sub>2</sub>

With the bare ansatz:

$$\lambda_\rho = [(-\beta_2 \log M)\rho + (\beta_3/\beta_2) \log \rho + C]^{-1} \quad (\text{III.4.15})$$

$$m_\rho = m \cdot \rho^{-\gamma} \quad a_\rho = 1 \quad (\text{III.4.16})$$

$$\beta_2 \equiv -2(N-1)/\pi \quad \beta_3 = 2(N-1)/\pi^2 \quad \gamma = (N-1/2)/(N-1) \quad (\text{III.4.17})$$

and  $C$  a large constant, the normalized Schwinger functions have a limit as  $\rho, \Lambda \rightarrow \infty$ , which correspond to finite values of  $\lambda_{ren}$ ,  $m_{ren}$  and  $a_{ren}$  in the BPHZ scheme.  $m_{ren}$  is close to  $m$ ,  $a_{ren}$  is close to 1, and  $\lambda_{ren}$  is small and non-zero, so that the theory is not trivial. Furthermore the corresponding theory is the Borel sum of its renormalized series (and therefore can be analytically continued to a disk  $C_R$ ; this requires to take  $C$  in (III.4.15) complex with a large real part). It is

therefore independent of the particular limiting process used to construct it, and in particular it obeys the O.S axioms for fermions.

We give only a sketch of the proof since it overlaps strongly with the previous constructions of section II.5 and III.3, and also because the reader has already the choice between several rather detailed published constructions [FMRS4] [GK4] [IM2]. Again we could avoid to consider running masses by performing full (useful and useless) mass renormalization. As discussed above, this requires an analysis at the level of one particle irreducible objects, rather than simply connected objects. We refer to [FMRS4] for such an approach.

We have to check first that the only counterterms of the theory are of the expected type, namely  $(\bar{\psi} \cdot \psi)^2$ ,  $\bar{\psi} \cdot \psi$  and  $\bar{\psi} \cdot (i \not{\partial})\psi$ . This is an analogue of Furry's theorem in electrodynamics, involving some algebraic manipulation of gamma matrices. For its proof we refer to [FMRS4, Lemma 2.1]. Remark that a four point function not of the form  $(\bar{\psi} \cdot \psi)^2$  is not identically 0, but it is convergent and does not need any renormalization.

As in the previous section, a phase space expansion is applied to (III.4.4). For each scale  $i$ , starting with  $i = \rho$  and ending at  $i = 0$  we perform first an horizontal cluster expansion with respect to  $C^i$  and the lattice  $\mathbf{D}^i$  of cubes of side size  $M^{-i}$ , then a fifth order vertical decoupling expansion, and finally a Mayer expansion with respect to the vacuum graphs and the internal domains of the two and four point functions, keeping their external cubes fixed.

At the end of the expansion we have explicit propagators or vertices (which realize the explicit connections of the expansion) and the remaining fields still have the structure of a determinant. This determinant is bounded using lemma III.4.1. Hence as announced there is no longer any distinct treatment of the “high” and “low momentum” fields.

We compute at each scale counterterms  $\delta\lambda_i$ ,  $\delta m_i$  and  $\delta a_i$  exactly like in the previous section. These counterterms are introduced in the vertical interpolations so as to renormalize all the local two and four point configurations generated by the expansion.

In the last slice we obtain a sum of usefully renormalized contributions with ordinary vertices and counterterm vertices; in particular there are 2-leg vertices with coefficients  $\sum_{j=i+1}^{\rho} \delta m_j$  and  $\sum_{j=i+1}^{\rho} \delta a_j$  respectively for a mass or wave function vertex with highest leg in slice  $i$ . To relate oneself to the BPHZ normalization conditions one has to resum all two point insertions into an effective propagator. This is possible (since the corresponding series is geometric). This effective propagator is of the form

$$e^{-V} [\not{p}(1 - \sum_{i=1}^{\rho} \delta a_i e^{-V}) + m_{\rho} - \sum_{i=1}^{\rho} \delta m_i e^{-V} + O(p^2)]^{-1} \quad (\text{III.4.18})$$

with  $V = (p^2 + m_{\rho}^2)$ .

The discrete flow of  $\lambda$ ,  $m$  and  $a$  is governed by the leading graphs, and to land on the desired renormalized coupling we must, like in section II.5, push the analysis up to third order graphs for the  $\beta$  function. These discrete equations for  $\lambda$  and  $a$  are:

$$\delta\lambda_i = -\beta_2 \log M \lambda_i^2 + \beta_2 \log M (\beta_2 \log M - \beta_3) + O\left(\frac{\log i}{i^4}\right) + \lambda_i^2 [O(e^{-(\rho-i)}) + O(e^{-i})] \quad (\text{III.4.19})$$

$$\delta a_i = \gamma_2 \log M \lambda_i^2 (1 + O(\lambda_i) + O(e^{-(\rho-i)})) + O(e^{-i}) \quad (\text{III.4.20})$$

where  $\gamma_2 = (2N - 1)/(2\pi)^2$  is a numerical constant corresponding to the graph  $B_0$  in Fig.II.5.1;  $\beta_3$  is again (as in the planar  $\phi_4^4$  theory, see lemma II.5.2) the sum of two contributions  $\gamma_3 = (N - 3/2)/\pi^2$  which corresponds to the analogue of the graph  $Q_3$  in Fig.II.5.1, and  $\delta_3 = 2\gamma_2$ , which corresponds to a reaction of the wave function renormalization on the coupling constant flow similar to the one due to the graph  $Q_5$  in Fig.II.5.1.

Only the discrete flow of the mass is significantly different from the “wrong sign” planar  $\phi_4^4$  flows. Indeed by parity considerations the mass renormalization is only logarithmically divergent. The leading contribution is given by a tadpole with mass counterterms insertion (Fig.III.4.1); the corresponding numerical factor is  $\gamma_1 = -(2N - 1)/\pi$ . If we used a resummed propagator with running mass  $m_i$ , the corresponding equation would be  $\delta m_i \simeq -\gamma_1 m_i \lambda_i \log M + O(\lambda_i^2)$ ; but since we leave the mass counterterms as interaction vertices we find an equivalent but slightly more complicated discrete flow equation, which is [IM2]:

$$\delta m_i = \left( \lambda_i \sum_{j=1}^{\rho} \delta m_j C(j-i) \right) [1 + O(\lambda_i) + O(e^{-(\rho-i)})] + O(e^{-i}) \quad (\text{III.4.21})$$

where

$$C(k) = \int \frac{d^2 p}{p^2} (e^{-(p^2 + m_\rho^2)} - 1) (e^{-M^{-2k}(p^2 + m_\rho^2)} - 1) \quad (\text{III.4.22})$$

In both points of view, if we define:

$$\gamma = -\gamma_1/\gamma_2 = \lim_{i \rightarrow \infty} \lim_{\rho \rightarrow \infty} \frac{\delta m_i}{m_i \lambda_i \beta_2} = (N - 1/2)/(N - 1) \quad (\text{III.4.23})$$

we have the approximate behavior  $\frac{\delta m_i}{m_i} \simeq -\frac{\gamma}{i}$ , which is consistent with the ansatz (III.4.16). More precisely we can show that equations (III.4.15-16) lead by an easy induction to the behavior:

$$\lambda_i = [(-\beta_2 \log M)i + (\beta_3/\beta_2) \log i + C + f(i)]^{-1} \quad (\text{III.4.24})$$

$$m_i = m \cdot i^{-\gamma} (1 + g(i)), \quad i \geq \text{const} \cdot C / (-\beta_2 \log M) \quad (\text{III.4.25a})$$

$$K^{-1} \leq m_i/m \leq K, \quad i \leq \text{const} \cdot C / (-\beta_2 \log M) \quad (\text{III.4.25b})$$

$$\left| \sum_{j=i+1}^{\rho} \delta a_j \right| \leq K \sup\{\lambda_i, M^{-2i}\}; \quad |a_j - 1| \leq 1/2 \quad \forall j \quad (\text{III.4.26})$$

for some large constant  $K$ , with  $|f(i)| \leq 1/2$  and  $|g(i)| \leq 1/2$ . All these behaviors assume that  $C$  is chosen sufficiently large in (III.4.15), and they are sufficient to verify the convergence of the usefully renormalized phase space expansion.

It remains to relate precisely the ultraviolet limit obtained in this way to the BPHZ conditions, which are:

$$\lambda_{\text{ren}} = S_4^A(C, m)(0, 0, 0, 0) \quad (\text{III.4.27})$$

$$(m_{\text{ren}})^{-1} = S_2^C(C, m)(0, 0) \quad (\text{III.4.28})$$

$$(a_{\text{ren}} = (m_{\text{ren}})^2 [-i \frac{d}{dp_0}] S_2^C(C, m)(p)|_{p=0}) \quad (\text{III.4.29})$$

One can check then that the parameters  $C$  and  $m$  in (III.4.15-16) correspond to an other renormalization scheme (which in terms of the Schwinger functions is implicit rather than explicit like BPHZ). More precisely, like in section II.5, we can study the map from  $(C, m)$  to  $(\lambda_{\text{ren}}, m_{\text{ren}})$ , and prove that for  $m$  close to  $m_{\text{ren}}$  and  $\text{Re}C$  large enough, we can obtain any prescribed value of  $(\lambda_{\text{ren}}, m_{\text{ren}})$  with  $\lambda_{\text{ren}}$  in a disk  $C_R$ . The resulting theory is analytic in this disk and again Taylor remainders at large order may be evaluated by the combination of the large order estimate of part II and the convergence of the phase space expansion, so that the theory is the Borel sum of its perturbative expansion.

This result is particularly welcome for this model, which is a full fledged field theory. Indeed it can be called upon to check that the ultraviolet limit constructed in this way is universal (independent of the technical details of the construction) and it allows a quick proof of the O.S. axioms: since different regularization schemes preserve different subsets of the axioms, it is easy to check the corresponding axioms. But by the Borel summability result, the theory constructed using these different regularizations is the same, hence verifies all axioms! For instance we can use the cutoffs (III.4.6) to check every axiom except reflection positivity (OS3 in section I.2), and a lattice cutoff to check this last axiom. Indeed lattice regularization preserves reflection positivity with respect to the symmetry hyperplanes of the lattice, which become denser and denser as the lattice spacing tends to zero; in the limit reflection positivity is recovered with respect to every hyperplane (here in dimension two these hyperplanes are simply straight lines).

## B Three dimensions

In the previous sections the technique of phase space expansion appears as a compromise between the perturbative expansion, which allows rather precise information on the model but diverges, and the functional integral which is a beautiful resummation of this expansion, but does not allow detailed estimates. However this main idea can be used of course in a more general context, and in particular phase space expansion can be used as a compromise between functional integration and an other expansion scheme than the ordinary perturbative one; what seems however necessary up to now for this method to work is the existence of a small parameter which, when tending to zero, drives the theory towards a gaussian one.

The Gross-Neveu model in three dimensions is an example of such a situation: it is no longer renormalizable in the ordinary perturbative sense but it is renormalizable in the sense of the  $1/N$  expansion, as is explained below. Furthermore its heuristic renormalization group analysis shows a nontrivial fixed point which is close to a gaussian one if  $N$ , the number of colors, is big enough. This fixed

point may be considered the remnant of the two dimensional asymptotic freedom of the model. Rigorous construction of the model using the phase space analysis is in progress [dCFdVMS], and we will discuss briefly here the content of the model, the expected results and the technical difficulties to overcome.

The starting point is the same lagrangian density than for the two dimensional model (III.4.1):

$$S(\bar{\psi}, \psi) = \bar{\psi} C^{-1} \psi - (\lambda/2N)(:\bar{\psi}\psi:)^2 \quad (\text{III.4.30})$$

with  $C^{-1} = (ia\partial + m)$ . The main change is that we are in a three dimensional space, so we can choose the Fermi fields to be 4 component spinors, and use the first three of the usual 4 by 4 gamma matrices of four dimensional space time. The usual  $\gamma_5$  matrix may be used also to introduce chirality in the model. An other change is that the parameter  $N$ , the number of color species, will be taken large, so that it is convenient to normalize the coupling constant as in (III.4.30).

One can no longer treat the model at a purely perturbative level like in the two dimensional case, because perturbation theory is not renormalizable in this case. To avoid an ill-defined fermionic integration one uses therefore the well known method of Matthews-Salam, in which the Fermi functional integral is exchanged for a bosonic integration over an auxiliary field  $\sigma$ . Formally one writes at each point:

$$e^{\frac{\lambda}{2N}(\bar{\psi}\psi)^2} = \int e^{-(1/2)\sigma^2 + \sqrt{\frac{\lambda}{N}}\sigma\bar{\psi}\psi} d\sigma \quad (\text{III.4.31})$$

and integrating over the Fermi fields, the partition function of the system becomes:

$$\begin{aligned} & \int e^{-\bar{\psi}(ia\partial + m)\psi - (\lambda/2N)(:\bar{\psi}\psi:)^2} \prod_x d\bar{\psi}(x)d\psi(x) \\ &= \int e^{-(1/2)\sigma^2 - \sqrt{\lambda N}TrC\sigma} \det(C^{-1} + \sqrt{\frac{\lambda}{N}}\sigma)^N \prod_x d\sigma(x) \end{aligned} \quad (\text{III.4.32})$$

(the partial Wick ordering is responsible for the extra linear term in  $\sigma$ ; without it the measure would blow up as  $N \rightarrow \infty$ ). This formula has now a natural gaussian limit as  $N \rightarrow \infty$ . Indeed we can write:

$$\begin{aligned} \det(C^{-1} + \sqrt{\frac{\lambda}{N}}\sigma)^N &= (\det C)^{-N} \det(1 + \sqrt{\frac{\lambda}{N}}C\sigma)^N \\ &\simeq e^{\sqrt{\lambda N}TrC\sigma - (\lambda/2)TrC\sigma C\sigma + O(1/\sqrt{N})} \end{aligned} \quad (\text{III.4.33})$$

Therefore up to a normalization the limit  $N \rightarrow \infty$  leads to the natural gaussian measure in  $\sigma$ :

$$\int e^{-(1/2)\sigma^2 - (\lambda/2)TrC\sigma C\sigma} \prod_x d\sigma(x) \quad (\text{III.4.34})$$

It is important to remark that this measure is no longer a simple local mass term in  $\sigma^2$  (a  $\delta$  function in  $x$  space); the non local correction is given by the bubble graph  $TrC(x-y)C(y-x)$ . To define more rigorously this  $N \rightarrow \infty$  limit we have to worry about ultraviolet finiteness and about the positivity of the corresponding gaussian measure, which is not obvious because the bubble graph is negative (in

field theory it is a well known fact that fermionic loops have a minus sign when compared to bosonic loops).

An ultraviolet cutoff is first applied to the fermionic propagator  $C$  which is therefore changed into  $C_\rho$  like in (III.4.6). The covariance of the measure (III.4.34) for the  $\sigma$  field (corresponding to the cutoff bubble graph) is:

$$\Gamma_\rho^0(p) = (1 + \lambda_\rho \pi_\rho(p))^{-1} \quad (\text{III.4.35})$$

with a well defined value:

$$\pi_\rho(p) = \frac{1}{(2\pi)^3} \int Tr C_\rho(k) C_\rho(p-k) d^3 k \quad (\text{III.4.36})$$

Now for a correct choice of  $\lambda_\rho$  as  $\rho \rightarrow \infty$  this covariance is positive so that the corresponding gaussian measure is well defined. In order to see this, it is convenient to rescale first the  $\sigma$  field by  $\sqrt{\lambda_\rho} \sigma \rightarrow \sigma$  so that  $\lambda_\rho$  disappears from the interaction. The  $\sigma$  covariance becomes  $\frac{1}{\lambda_\rho^{-1} + \pi_\rho(p)}$ , and the constant term  $\lambda_\rho^{-1}$  is chosen so as to control the zero momentum value of  $\pi_\rho$ , which is negative. More precisely if we define  $\lambda_\rho^{eff} = \frac{\lambda_\rho}{1 + \lambda_\rho \pi_\rho(0)}$ , and ask  $\lambda_\rho^{eff}$  to have a positive limit as  $\rho \rightarrow \infty$ , the covariance becomes

$$\frac{1}{(\lambda_\rho^{eff})^{-1} + \pi_{ren}(p)} \quad (\text{III.4.37})$$

The advantage is that  $\pi_{ren}(p)$  is now positive, hence for positive  $\lambda_\rho^{eff}$  the covariance is of positive type. It is easy to check that the corresponding asymptotic behavior of  $\lambda_\rho$  as  $\rho \rightarrow \infty$  is:

$$\lambda_\rho \simeq (c^{-1} + M^\rho \frac{e}{1 - M^{-1}})^{-1} \quad (\text{III.4.38})$$

In this way the limit  $N = \infty$  is a well defined gaussian theory in terms of the  $\sigma$  field.

This gaussian limit and the formula (III.4.33) is the correct starting point for the construction of the theory, and the expansion scheme will be based on the  $1/N$  rather than the perturbative expansion. To factorize the gaussian piece one uses the standard notation:

$$\det_{n+1}(1 + K) = \det(1 + K) e^{Tr(-K + \frac{K^2}{2} + \dots + (-1)^n \frac{K^n}{n!})} \quad (\text{III.4.39})$$

and one introduces a coupling constant in the form  $\lambda_\rho g_\rho^2$  where  $\lambda_\rho$  behaves as in (III.4.38) up to  $1/N$  corrections and  $g_\rho$  should converge to 1 as  $N \rightarrow \infty$ . Then a normalized Schwinger function, with test functions  $f_i, g_j$  is given by:

$$\begin{aligned} S_{\Lambda, \rho}(f_1, \dots, f_p; g_1, \dots, g_p) &= \frac{1}{Z} \int \bar{\psi}(f_1) \dots \bar{\psi}(f_p) \psi(g_1) \dots \psi(g_p) e^{\lambda_\rho g_\rho^2 \int_\Lambda (\bar{\psi} \psi)^2} d\mu_\rho(\bar{\psi}, \psi) \\ &= \frac{1}{Z'} \int \det(f_i \frac{1}{1 + K_\rho} g_j) \det_3(1 + K_\rho)^N d\mu_{\Gamma_\rho}(\sigma) \end{aligned} \quad (\text{III.4.40})$$

where  $K_\rho$  is the regularized kernel:

$$K_\rho(x, y) = \frac{g_\rho}{\sqrt{N}} C_\rho(x - y) \Lambda(y) \sigma(y) \quad (\text{III.4.41})$$

and

$$\Gamma_\rho(p) = (\lambda_\rho^{-1} + g_\rho^2 \pi_\rho(p))^{-1} \quad (\text{III.4.42})$$

The power counting using the covariance of the  $\sigma$  field is now better than for the initial perturbative expansion. Apart from vacuum diagrams, the only divergent diagrams are those with one  $\sigma$  external field (quadratically divergent), two Fermi or  $\sigma$  fields (linearly divergent) and two Fermi fields plus one  $\sigma$  field (logarithmically divergent). Diagrams with three external  $\sigma$  fields, although apparently divergent are in fact finite because of the vanishing of the trace of an odd number of  $\gamma$  matrices. These renormalizations basically correspond to the renormalization of the initial parameters (masses and coupling constant) of the model.

Therefore rewriting the theory in terms of the  $\sigma$  field and absorbing the first term in the  $\det_2$  in the definition of a gaussian measure for this field has changed the theory from non-renormalizable to just renormalizable. In terms of the initial perturbative series, this operation is simply the resummation of all chains of bubble diagrams (see Fig.II.5.4).

To proceed further towards the rigorous construction of the theory one has to introduce a momentum slicing, a phase space expansion and effective parameters for the theory (III.4.40), and to analyze the corresponding renormalization group flows. These flows are controlled as before by a few leading order graphs. The corresponding fixed point is not exactly gaussian but close to a gaussian for  $N$  sufficiently large, and the definition of the corresponding renormalized parameters is not particularly difficult; we refer the interested reader to [Sen]. Nevertheless the construction of the model is more difficult than for the previous examples because there is less positivity in terms of the bosonic  $\sigma$  field. This point is of great constructive importance and has no counterpart in the ordinary standard (perturbative) analysis, so let us try to sketch it briefly.

The interaction of the theory in the  $\sigma$  field point of view is the  $\det_3$ . We remember that for bosonic theory there is a stability problem at large fields, which show up in the definition of the functional integral and (in a phase space expansion) in the domination process. The interaction  $\det_3(1 + K_\rho)$  is much worse than a positive  $\phi^4$  interaction in this respect. Indeed the standard bound on a  $\det_3$  is:

$$\log(\det_3(1 + K)) \leq \frac{1}{2} \text{Tr}(K^2 + KK^*) \quad (\text{III.4.43})$$

$\text{Tr}K^2$  corresponds to the bubble graph (with two external  $\sigma$  fields) and is negative. But  $\text{Tr}KK^*$  is nothing but the bubble graph with both external fields at the same point (see Fig.III.4.2). Therefore we can rewrite (III.4.43) as:

$$|\det_3(1 + K)| \leq e^{\frac{1}{2}\pi_{ren}} \quad (\text{III.4.44})$$

where  $\pi_{ren}$  stands for the renormalized bubble graph. On the other hand the functional bound on the interaction is of the same order than the positivity of the gaussian measure (remember that in (III.4.42)  $\lambda_\rho$  is adjusted so as to cancel the zero momentum value of  $\pi(p)$  so that apart from a remaining constant, the positivity of the gaussian measure is exactly that of  $\pi_{ren}(p)$ ). Therefore the ultimate positivity of the theory is much more marginal than in previous models, and this leads to many technical complications.

In particular a new feature of the expansion is the decomposition of the functional integral according to whether the mean of the  $\sigma$  field is large or small. In the small field region, the  $\det_3$  can be expanded; in the large field regions the bound (III.4.44) is applied and compensated by the positivity of the gaussian; furthermore a small factor per cube is earned from the remaining piece of the gaussian (corresponding to the term  $\lambda_\rho^{eff}$  in (III.4.38)). The cluster expansions are not performed inside the large field regions, which can be treated as a single block.

However in addition there are also some important technical changes at the level of the interpolation parameters corresponding to the vertical and horizontal cluster expansions. In the  $\phi^4$  problem of section III.3 there was separate positivity of the gaussian measure and of the interaction, so it was enough to write independent interpolations on both of them which preserved separately the positivity of each of them. Now by (III.4.34) it is only the combination of the  $\det_3$  interaction and of the gaussian measure which is positive hence the horizontal and vertical expansions cannot be fully independent, but must be defined in such a way that for any value of the interpolating parameters the combined interaction and gaussian measure remains positive. This is possible, but requires an inductive definition of these parameters, which takes into account already derived connections. We remark that like for the horizontal tree like expansion, it seems a general property that inductive interpolations which are minimal in the sense that they never build redundant connections are also optimal from the point of view of preserving positivity requirements.

This ends our brief description of the Gross-Neveu model in three dimensions, for which a detailed construction should be soon available [dCFdVMS].

### III.5 The ultraviolet problem in non-abelian gauge theories

*I see nothing wrong with it because any nontrivial idea  
is in a certain sense correct.*

– A. M. Polyakov, Gauge Fields and Strings.

#### A) Introduction

The efforts to understand better renormalization theory should culminate in a rigorous solution of the ultraviolet problem for non-abelian gauge theories. Most physicists are convinced that the problem is well understood and void of any surprises, because of its asymptotically free character. However there is only one rigorous program of study of this problem completed so far, the one of Balaban [Ba2-9]. This program defines a sequence of block-spin transformations for the pure Yang-Mills theory in a finite volume on the lattice and shows that as the lattice spacing tends to 0 and these transformations are iterated many times, the resulting effective action on the unit lattice remains bounded. From this result the existence of an ultraviolet limit for *gauge invariant* observables such as “smoothed Wilson loops” should follow, at least through a compacity argument using a subsequence of approximations; but the limit is not necessarily unique. Clearly this is a point which requires further work.

Although very impressive, Balaban’s work reaches the limits of human communicability, partly because the use of the lattice regularization is the source of many technical complications and partly because the results are scattered over many publications. Also it does not address the problem of constructing the expectations values of products of the field operators in a particular gauge (the Schwinger functions), because these are not gauge invariant observables. These remarks also apply to the related program of Federbush [Fed2-7] (an other impressive task which, like European common market or the Channel tunnel, is scheduled for completion by 1993). It is true that physical quantities should be gauge invariant. Nevertheless the gauge fixed framework is obviously the most convenient for perturbative computations, and one can consider in fact that the ultraviolet problem for the Yang-Mills *field* theory is not yet understood until this point is clarified.

In collaboration with J. Feldman, J. Magnen and R. Sénéor, we tried also our own study of this problem with the phase space method described in this book. Our ambitious and perhaps naive goal was to construct the Schwinger functions of the field e.g. in the Feynman or Landau gauge, with an infrared cutoff. In spite of hard work, at least in term of the hours spent, we did not succeed. The functional integrals obtained always lacked sufficient positivity for control. For long we hoped that at some point in the construction the phase space analysis would “kill” the related Gribov problem, or “chop” it into manageable subpieces, but in the author’s present opinion this hope was unfortunately ill-founded.

We do not have at the moment many intermediate results of such obvious value that we owe to explain them to the community. Nevertheless at the perturbative level we are convinced that our approach, which is based on a gauge breaking

regularization, can be used to recover the results on renormalizability of the Yang-Mills perturbation theory, without the use of dimensional regularization (which is not a useful constructive device up to now), and with explicit uniform bounds at large order such as those of part II. Furthermore with this method, one can rewrite the bare or renormalized series as an effective series with useful renormalization at the level of power series in the renormalized coupling like in section II.4. At this level of formal power series the running coupling constant seems clearly to display asymptotic freedom.

At the constructive level, the bare theory with cutoff that we consider is well defined (and this is somewhat non-trivial, since it depends on the shape of the ultraviolet cutoff that we use). But whether the ultraviolet limit exists remains a widely open question. Also if it exists it is not clear that it obeys the standard (perturbative) renormalization group computations: in [DeZw] it is claimed that for the theory in the Landau gauge this is not the case at least for the two point function. Clearly this surprising claim deserves better attention. In order to reconcile it with Balaban's work, and the standard "Monte-Carlo" wisdom, one could perhaps imagine that only gauge invariant observables behave according to the standard renormalization group.

In conclusion it is the goal of this section to explain our tentative approach to the Yang-Mills ultraviolet problem, and to summarize the little we know about the Gribov problem, in relation to our point of view.

## B) The model

We consider the pure Yang-Mills theory with an infrared cutoff, which we never try to lift. This cutoff may be imposed on the propagator, or we could consider the theory on a finite volume with some boundary conditions, or on the sphere  $S^4$ , the torus  $\Lambda = \mathbb{R}^4/\mathbb{Z}^4$  or an other compact Riemannian four-dimensional manifold. Naive infrared regularization breaks gauge invariance, but compactification of space and the choice of a particular principal bundle with fiber  $G$  defines an unbroken group of gauge transformations. For instance in the case of the torus with the trivial  $SU(2)$  bundle, the gauge transformation are simply the functions  $x \rightarrow g(x)$  from  $\mathbb{R}^4$  to  $G$  which are periodic with period lattice  $\mathbb{Z}^4$ . The momentum space corresponds to discrete Fourier analysis on the dual lattice  $\Lambda^* = \mathbb{Z}^4$ . Moreover the constant fields or the zero mode in Fourier space is deleted in all our functional integrals, hence there is no infrared problem.

For the pure  $SU(2)$  Yang-Mills theory the vector potential is a field  $A_\mu^a$ ,  $\mu = 1, \dots, 4$ ,  $a = 1, 2, 3$  with Lorentz (greek) indices and Lie algebra (latin) indices (the group is noted  $SU(2)$  and the algebra  $su(2)$ ). Geometrically  $A$  is a connection on the considered principal bundle; again in the case of the trivial  $SU(2)$  bundle one can consider that each  $A_\mu$  is simply a function with values in  $su(2)$ . Our conventions are those of [IZ], which we recall briefly; later to simplify the notations we will forget indices most of the time. We write  $A = \sum_{a=1}^3 A^a t_a$ , with  $t_a = (i\sigma_a/2)$  where the  $\sigma$ 's are the three usual hermitian Pauli matrices. With this convention the covariant derivative is  $D_\mu = \partial_\mu - \lambda[A_\mu, .]$ . We have  $\text{Tr } t_a t_b = -\frac{\delta_{ab}}{2}$ . The field curvature is:

$$F_{\mu\nu} = (\partial_\mu A_\nu - \partial_\nu A_\mu) - \lambda[A_\mu, A_\nu] = (\partial \wedge A - \lambda[A, A]) \quad (\text{III.5.1})$$

$\lambda$  being the coupling constant; the second notation is a condensed one in which indices are omitted (and  $\partial\wedge$  is the exterior derivative). Remark that in the three dimensional su(2) space, the commutator is a wedge product:  $[A_\mu^a, A_\nu^b] = \epsilon_{ab}^c A_\mu^a A_\nu^b$ . The pure Yang-Mills action is (for Euclidean canonical metric on the flat torus the raising of "Lorentz" indices is trivial so that  $F_{\mu\nu} = F^{\mu\nu}$ ):

$$-\frac{1}{2} \int_{\Lambda} d^4x Tr F_{\mu\nu} F^{\mu\nu} = \frac{1}{4} \int_{\Lambda} d^4x \sum_a F_{\mu\nu}^a F^{\mu\nu a} \quad (\text{III.5.2})$$

To simplify, we define a scalar product  $\langle A, B \rangle$  on space time tensors of the same type with values in the Lie algebra  $A$  and  $B$ , by the convention that a trace is taken over all correspondent space time indices and minus a trace over group indices, so that it is positive definite with a factor 1/2 in component notation. We also write simply  $A^2$  for  $\langle A, A \rangle$ , and with this convention we can write the action as  $\frac{1}{2} \int_{\Lambda} F^2$ . We distinguish between the quadratic, trilinear and quartic pieces of  $F^2$ , writing:

$$F^2 = F_2 + \lambda F_3 + \lambda^2 F_4 \quad (\text{III.5.3})$$

This action is invariant under the gauge transformations:

$$A \rightarrow A^g ; (A^g)_\mu = g A_\mu g^{-1} + \partial_\mu g \cdot g^{-1} \quad (\text{III.5.4})$$

In what follows these gauge transformations are limited to a particular topological sector, for instance the functions from the compact space to  $G$ . It is often useful to consider the infinitesimal gauge transformations  $\epsilon$  with values in the Lie algebra, which are tangent to the gauge transformations; the corresponding formula is:

$$A \rightarrow A^\epsilon ; (A^\epsilon)_\mu = A_\mu + D_\mu \epsilon \quad (\text{III.5.5})$$

where  $D = \partial - \lambda [A, .]$  is the covariant derivative. Our starting point is the Yang-Mills theory in the Feynman gauge. In this gauge there is an additional factor  $e^{-\langle \partial \cdot A, \partial \cdot A \rangle}$  which is the gauge fixing term, and a Faddeev-Popov determinant which is  $\det \partial \cdot D$ . This determinant can formally be written in terms of ghosts fields  $\eta$  and  $\bar{\eta}$  as

$$e^{\partial_\mu \bar{\eta}_a (\partial_\mu \eta_a - \lambda \epsilon_{abc} A_{\mu,b} \eta_c)} \quad (\text{III.5.6})$$

The propagator for the gauge and ghosts fields are respectively

$$\frac{\delta_{\mu\nu} \delta_{ab}}{p^2} , \quad \frac{\delta_{ab}}{p^2} \quad (\text{III.5.7})$$

and the Feynman rules are exactly the standard ones that can be found in [IZ] (after rotation to Euclidean space).

The class of ultraviolet cutoffs we consider is defined as follows.  $\kappa$  is a fixed function which is 1 near 0 and decreases at infinity. For instance it could be an exponentially decreasing function or a  $C_0^\infty$  function, which is 0 for  $|p| \geq 2$  and is 1 for  $|p| \leq 1$  (the  $C_0^\infty$  character is perhaps not essential but it should be such that the slices built out of it (see (III.5.10) have good spatial decay and it might be interesting to have also good momentum conserving properties).

Then we define our scaled momentum cutoff  $\kappa_\rho$  to be:

$$\kappa_\rho(p) = (1/2)[\kappa(\alpha p M^{-\rho}) + \kappa(p M^{-\rho})] \quad (\text{III.5.8})$$

where  $\alpha$  is a small constant (this unusual form leads to a stabilizing  $A^4$  counterterm, as shown below). We can consider the well defined normalized gaussian measure in  $A$ ,  $d\mu_\rho$ , whose propagator is:

$$C_\rho = \delta_{\mu\nu} \delta_{ab} \frac{\kappa_\rho}{p^2} \quad (\text{III.5.9})$$

and create the first momentum slice by writing, as before:

$$C^\rho = C_\rho - C_{\rho-1} \quad (\text{III.5.10})$$

The same cutoff and slicing is used for the ghost propagator. We will use the notation

$$\det_\rho(A) \equiv e^{-\lambda \partial_\mu \bar{\eta}_a \epsilon_{abc} A_{\mu,b} \eta_c} d\nu_\rho(\bar{\eta}, \eta) \quad (\text{III.5.11})$$

where  $d\nu_\rho$  is the formal free measure for the anticommuting fields  $\bar{\eta}$  and  $\eta$  with propagator  $\delta_{ab} \frac{\kappa_\rho}{p^2}$ . This notation is useful, but the true rigorous definition of the corresponding regularized determinant  $\det_\rho$  is in fact as in the previous section on the Gross-Neveu model, through convergent power series. There is indeed both an ultraviolet and an infrared cutoff, so by an analogue of Lemma III.4.1, for any smooth sample field  $A$  the power series for the determinant  $\det_\rho(A)$  indeed has infinite radius of convergence.

Corresponding to the slicing (III.5.10) there is an orthogonal decomposition of the field  $A^\rho$  which is the random variable associated to  $d\mu_\rho$  as  $A_\rho = A^\rho + A_{\rho-1}$  and of the ghosts fields  $\bar{\eta}_\rho = \bar{\eta}^\rho + \bar{\eta}_{\rho-1}$ ,  $\eta_\rho = \eta^\rho + \eta_{\rho-1}$ , but this is again only a convenient notation for the manipulation of determinants.

From standard renormalization group analysis we learn that in order to get a finite non trivial renormalized theory at the unit scale of our finite box, we should use a bare coupling constant which has the usual asymptotic behavior with  $\rho$  implied by asymptotic freedom. Hence a good ansatz for the bare coupling  $\lambda_\rho$  should be:

$$\lambda_\rho = \frac{1}{\beta_2(\text{Log}M)\rho + \beta_3/\beta_2 \log \rho + C} \quad (\text{III.5.12})$$

where  $C$  is a large constant, and  $\beta_2$  and  $\beta_3$  are the usual first non vanishing coefficients of the  $\beta$  function, whose numerical value is given in standard textbooks like [IZ]. Then one hopes that the renormalized coupling constant  $\lambda_{ren}$ , which should be defined as the last one in a sequence of effective constants, is finite and arbitrarily small as  $C$  becomes arbitrarily large by the same mechanism than in section II.5 (if perturbative renormalization group analysis turns out to be correct).

One could believe that the bare theory to start with is simply obtained by multiplying the well defined measure  $\det_\rho d\mu_\rho$  by the interaction terms  $e^{(\lambda F_3 + \lambda^2 F_4)}$ . However first it would not be obviously well defined since even with a cutoff the cubic term  $F_3$  is not positive and to bound it would require extracting the full positivity of the gaussian measure. Second, the true bare theory must be more complicated to respect gauge invariance in the limit  $\rho \rightarrow \infty$ . Indeed our ultraviolet

cutoff breaks gauge invariance. These gauge breakings must be compensated by appropriate gauge-variant counter terms. In fact only the relevant and marginal counterterms must be included because they are the only ones whose effect on finite scales does not vanish as  $\rho \rightarrow \infty$ . We show now how to compute these effects.

### C) Computation of the counterterms due to the ultraviolet cutoff

Our ultra violet cutoff does not break global SU(2) or Euclidean invariance (small Euclidean breaking effects nevertheless occur due to the infrared cutoff; for instance in the case of a torus there exist such effects due to the lattice structure of  $\Lambda^*$ , but they are tied to the unit scale and do not need counterterms). Therefore the only new relevant or marginal operators that we should consider are  $-\text{Tr}A_\mu A_\mu$ ,  $(-\text{Tr}A_\mu A_\mu)^2$ ,  $(-\text{Tr}A_\mu(-\Delta)A_\mu)$  and  $-\text{Tr}(\partial_\mu A_\mu)^2$  which we abbreviate respectively as  $A^2$ ,  $A^4$ ,  $A(-\Delta)A$  and  $(\partial A)^2$  (recall the convention that traces are definite negative). This is only true for SU(2) theory; for an SU(N) theory there would be a longer list of operators to consider and the analysis would be more complicated.

In fact our gauge breaking cutoff also disturbs the magic relation  $Z_2 Z_4 = Z_3^2$  which relates the multiplicative renormalization of  $F_2$ ,  $F_3$  and  $F_4$  in  $F^2$  and expresses the fact that up to a rescaling of  $A$  only the coupling constant  $\lambda$  is renormalized [IZ]. To correct this problem, using the possibility of rescaling  $A$ , we need only to introduce a single counterterm, for instance of the type  $F_4$ .

Therefore the counterterms that we introduce are:

$$e^{CT} = e^{-a_\rho \int_\Lambda A^4 - b_\rho \int_\Lambda A^2 - c_\rho \int_\Lambda A(-\Delta)A - d_\rho \int_\Lambda (\partial A)^2 - e_\rho \int_\Lambda F^4} \quad (\text{III.5.13})$$

The relevant counterterm  $b_\rho \int_\Lambda A^2$  must be fine tuned exactly to have a renormalized mass which is zero. This is the same problem than fixing the critical bare mass in infrared  $\phi_4^4$  and should be solved either by a fixed point argument as in section III.3 above or using a full renormalization of the two point function (and a one particle irreducible analysis) as in [FMRS5]. For the marginal counterterms, an analysis to lowest order in perturbation theory is in fact enough for our purpose (because of asymptotic freedom, further orders again should give no contributions to finite scales in the limit  $\rho \rightarrow \infty$ ). We obtain:

#### Lemma III.5.1

$$a_\rho \simeq a\lambda_\rho^4, \quad b_\rho \simeq bM^{2\rho}\lambda_\rho^2, \quad c_\rho \simeq c\lambda_\rho^2, \quad d_\rho \simeq d\lambda_\rho^2, \quad e_\rho \simeq e\lambda_\rho^4. \quad (\text{III.5.14})$$

Furthermore by choosing the cutoff of the form (III.5.8) with  $\alpha$  large enough (depending on the shape of  $\kappa$ ), the coefficient  $a$  is strictly positive\*.

**Proof** We recall the Feynman rules for the pure SU(2) gauge theory in the Feynman gauge [IZ]. The propagators are given in (III.5.7). The interaction vertices are

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\* It is not clear whether a cutoff for which  $a$  would be negative (or zero) can be used in a constructive way. The answer may depend on adding irrelevant counterterms of higher order generated by the cutoff, which may stabilize the theory. The analysis of globally invariant such terms becomes more and more complicated as the order increases and we will therefore not try to explore this possibility here.

of three kinds. For simplicity we always forget to write the overall multiplication factor (of  $2\pi$ ) and the  $\delta$  function which expresses momentum conservation which equips them. These three kinds of vertices are then pictured in Fig.III.5.1.

We concentrate on the computation of the  $A^4$  counterterm, which is the most interesting, and include also the computation of the  $A^2$  counterterm. The other ones are less interesting and left to the reader.

At one loop, which also means at order  $g^4$  in perturbation theory, there are 4 graphs which may contribute to the  $A^4$  term. They are pictured in Fig III.5.2 and called  $G_1$ ,  $G_2$ ,  $G_3$  and  $G_4$ . To compute their contribution, we may assume by symmetry that in all four external legs, both the space time and group indices are equal to 1.

### a) Computation of $G_1$

The graph is obtained by applying 4 derivatives  $\frac{\partial}{\partial A_1^1}$  on  $(1/2!)(-F^2/4)^2$ . The result is  $3(\partial^2 F^2/4)^2$  where derivatives are taken with respect to  $A_1^1$ . The only non vanishing pieces come from the derivatives acting on the commutator in  $F$ , hence  $\partial^2 F^2/4$  gives  $(1/2)(\partial F)^2$ . Moreover we have  $\partial F_{\alpha\beta}^c = \epsilon^{c1b}[A_{\beta\delta}^b\delta_{\alpha 1} - A_{\alpha\delta}^b\delta_{\beta 1}]$ , where  $\epsilon$  is the usual antisymmetric tensor. But remark that if  $\alpha = \beta = 1$  the term vanishes. Hence when developing the square  $(1/2)(\partial F)^2$  the cross terms vanish. Therefore this square gives  $(\epsilon^{c1b})^2(A_\beta^b)^2\delta_{\alpha 1}$ ,  $\beta \neq 1$ . There are now two possible Wick contractions, a sum over three values (2,3 and 4) for  $\beta$  and a sum over 2 values (2 and 3) for  $b$ . Collecting all factors we obtain a positive coefficient  $3.2.3.2 = 36$  in front of the integration over the loop momentum of the two propagators of  $G_1$ .

### b) Computation of $G_2$ .

We apply 4 derivatives on  $(1/3!)(-F^2/4)^3$ . The result is  $-6(\partial^2 F^2/4)(\partial F^2/4)^2$  where derivatives are again with respect to  $A_1^1$ . The term in  $\partial^2 F^2/4$  is the same as before, hence gives  $(\epsilon^{c1b})^2(A_\beta^b)^2\delta_{\alpha 1}$ ,  $\beta \neq 1$ . But we have now two trilinear vertices in  $\partial F^2/4$  hence terms with derivative couplings; remark that a partial derivative  $\partial_\mu$  can be replaced by  $-ik_\mu$ . The computation of this term leads to two identical vertices, one which gives  $\epsilon^{1mn}A_\mu^n[\partial_1 A_\mu^m - \partial_\mu A_1^m]$ , and the other with  $m, n, \mu$  respectively replaced by  $p, q, \lambda$ . In the Wick contraction schemes we can first contract to form the line between these two trilinear vertices. Since the two half legs of the remaining vertex bear the same index  $\beta \neq 1$ , a tedious computation gives that the only term compatible with future contractions is  $(\epsilon^{1mn})^2(A_\mu^m)^2[4k_1^2 + k_\mu^2]$ . Using Euclidean symmetry, this is equivalent to  $(\epsilon^{1mn})^2(A_\mu^m)^2[5k_1^2]$ . Contracting with the remaining vertex, we have now as before two possible Wick contractions, a sum over three values (2,3 and 4) for  $\beta$  and a sum over 2 values (2 and 3) for  $b$ . Collecting all factors we obtain a negative coefficient  $-6 \cdot 2 \cdot 3 \cdot 2 \cdot 5 \cdot k_1^2 = -90 \cdot 4k_1^2$ , again equivalent by Euclidean symmetry to  $90k^2$  in front of the integration over the loop momentum of  $G_2$ .

### c) Computation of $G_3$

We apply 4 derivatives on  $(1/4!)(-F^2/4)^4$ . The result is  $+(\partial F^2/4)^4$  where derivatives are again with respect to  $A_1^1$ . The term in  $\partial F^2/4$  gives the same trilinear vertex as before, hence gives  $\epsilon^{1mn}A_\mu^n[\partial_1 A_\mu^m - \partial_\mu A_1^m]$ . In the Wick contraction schemes we can first choose one particular leg of vertex 1 to form a first line between two trilinear vertices. To choose the vertex (2,3 or 4) to which this leg contracts gives a factor 3. After this contraction has been performed, the

line equipped with two not yet contracted fields gives a term  $(\epsilon^{1mn})^2[2k_1^2(A_\mu^m)^2 + k_\mu^2(A_1^m)^2 - 3k_1k_\mu A_1^m A_\mu^m]$ . Here we can assume  $\mu \neq 1$ . We can now contract once more to create one line between the two remaining vertices, and this can be done in all possible ways, hence gives a different term, which is  $(\epsilon^{1mn})^2[4k_1^2(A_\mu^m)^2 + k_\mu^2(A_1^m)^2 - 6k_1k_\mu A_1^m A_\mu^m + k_\mu k_\lambda A_\mu^m A_\lambda^m]$ . We can assume that  $\mu \neq 1$  in the first three terms and that  $\mu = \lambda = 1$  is excluded in the last one. It remains to contract together both expressions. We have as before two possible Wick contractions, a sum over three values (2,3 and 4) for  $\mu$  and a sum over 2 values (2 and 3) for  $m$ . After collecting all factors, taking into account Euclidean symmetry we obtain a positive contribution  $9 \cdot 12 \cdot k_1^4 + 10 \cdot 12k_1^2k_2^2$  in front of the integration over the loop momentum of  $G_3$ . Converting it into units of  $(k^2)^2$ , we find a final combinatoric factor  $9 \cdot 12/8 + 10 \cdot 12/24 = 55.5$ .

#### d) Computation of $G_4$

We apply 4 derivatives on  $(1/4!)(F.P.)^4$ , where  $F.P.$  means the Faddeev-Popov term  $\partial_\mu \bar{\eta}_a (D_\mu \eta)_a$ , with  $D$  the covariant derivative. The result is  $(\partial_1 \bar{\eta}_a \epsilon_{ab1} \eta_b)^4$ . The combinatoric is easier. We obtain a factor 6 for the Wick contractions, a factor 2 for summations over latin indices and a minus sign corresponding to the fermionic loop, which comes from reordering correctly the anticommuting fields  $\eta$  and  $\bar{\eta}$ . Hence the contribution is  $-12 \cdot k_1^4$  in front of the integration over the loop momentum of  $G_4$ . Applying the same conversion rate, we obtain in units of  $(k^2)^2$  a final combinatoric factor of  $-1.5$ .

Remark that all 4 coefficients add up to 0. This is a particular case of the famous miracle of renormalizability (at one loop...) of four dimensional gauge theories.

Let us perform now a similar analysis for the  $A^2$  counterterm. There are three graphs contributing at order  $g^2$ , pictured in Fig III.5.3.

The first graph,  $G'_1$ , gives  $\partial\partial F^2/4 = (1/2)(\partial F)^2 = -(1/2)(\epsilon^{a1c}[A_\lambda^c \delta_{\mu1} - A_\mu^c \delta_{\lambda1}])^2$  which is non zero only for  $\mu\nu \neq 1$ . The contribution is -6 in front of the integration over the loop momentum. The second graph,  $G'_2$ , gives  $\partial\partial(1/2)(F.P.)^2 = (\partial F.P.)^2 = (\partial_1 \bar{\eta}_a \epsilon_{ab1} \eta_b)^2$ . The contribution is positive because the minus due to the fermion loop is compensated by the  $i^2 = -1$  coming from the derivatives on the fields. The contribution is  $2k_1^2 = (1/2)k^2$  in front of the integration over the loop momentum. The last graph,  $G'_3$ , is given by  $\partial\partial(1/2)(F^2/4)^2 = (\epsilon^{1mn} A_\mu^n [\partial_1 A_\mu^m - \partial_\mu A_1^m])^2$  (which is non zero only for  $\mu \neq 1$ ). The contribution is  $18k_1^2 = (9/2)k^2$  in front of the integration over the loop momentum.

The result for the  $A^2$  term is then  $(-6+1/2+9/2)=-1$  times the loop integration.

To complete the Lemma, we want to study the sign of the  $A^4$  counterterm. Let us explain why it is important to us. Our strategy is to cancel explicitly the  $A^4$  and  $A^2$  contributions due to the gauge breaking character of our ultraviolet cutoff by appropriate counterterms. Remark that strictly speaking, only the  $A^2$  contribution diverges as  $\rho \rightarrow \infty$  and requires a counterterm (for the  $A^4$  term the coefficient of the divergent piece is 0, as computed above). However this  $A^2$  counterterm is positive (since the contribution is negative, see the -1 above). This is dangerous for stability estimates. We will use the (finite)  $A^4$  counterterm to control this dangerous  $A^2$  term and stabilize the theory. But this requires that we use an ultraviolet cutoff such that the  $A^4$  counterterm is negative, hence such that

the total  $A^4$  contribution induced by the cutoff is positive. As a consequence of our expansion the leading contribution is the one-loop contribution; we want its sign to be positive. We show now that this is possible provided the covariance defining our ultraviolet cutoff is taken to be the arithmetic mean of two covariances, one with cutoff  $K \cdot M^{2\rho}$  and the other with cutoff  $(1/K) \cdot M^{2\rho}$ ,  $K$  being a large constant. The case (III.5.8) is then similar, with  $\alpha$  playing the rôle of  $K^2$ , up to a redefinition of the unit scale.

Let  $\kappa(p)$  be the ultraviolet cutoff function in momentum space. The phenomenon that we will discuss here is universal, and does not depend of the particular form of  $\kappa$ . Using the coefficients computed in the preceding section, the one loop contribution to the  $A^4$  term is, for a single cutoff  $\kappa_\rho(p) = \kappa(pM^{-\rho})$  (all our integrals are infrared regularized and "finite" means finite as  $\rho \rightarrow \infty$ ):

$$\int \frac{d^4 p}{p^4} 36\eta^2(pM^{-\rho}) - 90\eta^3(pM^{-\rho}) + 54\eta^4(pM^{-\rho}) = 0 \cdot \rho + \text{finite terms} \quad (\text{III.5.15})$$

For the arithmetic mean of two cutoffs, the corresponding contribution is:

$$\begin{aligned} & \int d^4 p \left[ \frac{36}{4} (\kappa(KpM^{-\rho}) + \kappa(K^{-1}pM^{-\rho}))^2 \right. \\ & \left. - \frac{90}{8} (\kappa(KpM^{-\rho}) + \kappa(K^{-1}pM^{-\rho}))^3 + \frac{54}{16} (\kappa(KpM^{-\rho}) + \kappa(K^{-1}pM^{-\rho}))^4 \right] \end{aligned} \quad (\text{III.5.16})$$

But we have the following rules for any reasonable cutoff (and certainly for the ones we use), and for integers  $q$  and  $r$ :

$$\int \frac{d^4 p}{p^4} (\kappa^q(K \cdot pM^{-\rho})) = -\log KM^{-\rho} + \text{finite terms if } q \geq 1 \quad (\text{III.5.17})$$

and, if  $K \gg 1, M > 1$ :

$$\begin{aligned} & \int \frac{d^4 p}{p^4} (\kappa^q(K \cdot pM^{-\rho})(\kappa^r(K^{-1} \cdot pM^{-\rho})) = -\log KM^{-\rho} + \text{finite terms if } q \geq 1 \\ & \int \frac{d^4 p}{p^4} (\kappa^q(K \cdot pM^{-\rho})(\kappa^r(K^{-1} \cdot pM^{-\rho})) = -\log K^{-1}M^{-\rho} + \text{finite terms if } q = 0, r \geq 1 \end{aligned} \quad (\text{III.5.18})$$

As a consequence the one loop  $A^4$  contribution behaves as

$$\left( -\frac{36 \times 2}{4} + \frac{90 \times 6}{8} - \frac{54 \times 14}{16} \right) \log K + \text{finite terms} = 2.25 \log K + \text{finite terms} \quad (\text{III.5.19})$$

where "finite terms" now means terms which are uniformly bounded both as  $\rho$  and  $K$  tends to  $+\infty$ . Therefore taking  $K$  large enough (depending on the details of our cutoff, which are responsible for the particular value of the finite terms) we can always achieve our goal of a positive total  $A^4$  contribution, hence of a negative stabilizing counterterm.

In conclusion our starting ansatz for the bare theory is:

$$e^{J_\rho} d\mu_\rho(A) d\nu_\rho(\bar{\eta}, \eta) \quad (\text{III.5.20})$$

with

$$J_\rho = J(A_\rho, \bar{\eta}_\rho, \eta_\rho) = -\frac{1}{2}(\lambda_\rho \int_\Lambda F_3 + \lambda_\rho^2 \int_\Lambda F_4) - \lambda_\rho \partial_\mu \bar{\eta}_a \epsilon_{abc} A_{\mu,b} \eta_c + CT \quad (\text{III.5.21})$$

with  $\lambda_\rho$  as in (III.5.12). This starting point is now clearly well defined because we have both finite volume and ultraviolet cutoff. Hence the sample fields are smooth. At large field the exponential of the interaction is clearly bounded for any sample field, in every direction of the configuration space. This is true because the leading terms at large field are the  $A^4$  and  $F_4$  terms which, thanks to the sign of the  $A^4$  term, are respectively positive definite and positive (the determinant  $\det_\rho(A)$  may be bounded at large  $A$  for instance by an analogue of (III.4.43), hence by a quadratic term). Remark however that it is only for fields of order  $\lambda^{-1}$  that the  $A^4$  term provides convergence, so this term does not confine the field in the true perturbative region ( $A \ll \lambda^{-1}$ ).

#### D) Perturbative results

For perturbative results, we can expand the  $e^{J_\rho}$  term in (III.5.20) as a formal (bare) power series, which integrated with the gaussian measure  $d\mu_\rho d\nu_\rho$  gives a bare perturbative expansion. Then we must rewrite this bare expansion as an effective expansion following the method of section II.4. However perturbative results are limited to statements about formal power series, and we cannot in this context justify asymptotic freedom. In particular we must abandon the ansatz (III.5.12) and consider the coupling  $\lambda_\rho$  simply as the parameter of the series. The violations of gauge invariance due to the ultraviolet cutoff are still of the form (III.5.13) but we can no longer use asymptotic freedom (III.5.12) and Lemma III.5.1 to argue that for instance only the one loop contribution has to be included in the  $A^4$  counterterm. For perturbative results we need to express  $a_\rho$ ,  $b_\rho$ ,  $c_\rho$ ,  $d_\rho$  and  $e_\rho$  as full formal power series in  $\lambda_\rho$ , computing each order as is done explicitly above for the one loop order. However there is an advantage: we have no longer positivity requirements for the coefficient  $a$  in (III.5.14) hence we are no longer limited to cutoffs of the type (III.5.8). We can use any shape we want for the ultraviolet cutoff function.

If we compare to the case of  $\phi_4^4$  treated in section II.4, there is one main additional problem. We want to show that the resulting usefully renormalized theory has not as many running parameters as the naive number of relevant or marginal operators (which, including non-Euclidean invariant ones, is in the order of the hundreds even for pure SU(2), since the field has twelve components). First by Euclidean and global SU(2) invariance of the slice cutoffs we show that there is no divergence associated to subgraphs whose external leg indices correspond to non-Euclidean or non-global SU(2) invariant operators. Therefore the corresponding subgraphs do not require renormalization at all, and in particular no useful renormalization; no running parameters are therefore associated to the corresponding operators.

Still among the coefficients of the remaining operators in the field  $A$ , which are  $A^2$ ,  $A^4$ ,  $(\partial A)^2$ ,  $F_2$ ,  $F_3$ , and  $F_4$  ( $A(-\Delta)A$  being a combination of  $F_2$  and  $(\partial A)^2$ ) we need some relations to hold in order to show that the effective parameters at scale  $i$  reconstruct a theory similar to (III.5.20) but with running parameters. There are also relations to be checked involving the divergent operators with ghosts; the

ghost wave function renormalization, which we abbreviate as  $\bar{\eta}\partial^2\eta$  and the ghost-field coupling in (III.5.11), which we abbreviate as  $\partial\bar{\eta}A\eta$  (other ghost terms, in particular the relevant ghost mass renormalization are zero by unbroken symmetries).

Since a wave function renormalization is allowed both for the field and the ghosts, there are two main relations to be checked, which up to a rescaling of the fields express that only  $\lambda$ , the coupling constant, is renormalized; these are the two relations (12-126a) of [IZ], which we prefer to rewrite (with hopefully slightly more logical notations) as:

$$\frac{Z_4}{Z_3} = \frac{Z_3}{Z_2} = \frac{\tilde{Z}_3}{\tilde{Z}_2} \quad (\text{III.5.22})$$

where  $Z_2$ ,  $Z_3$  and  $Z_4$  are the multiplicative coefficients of  $F_2$ ,  $F_3$  and  $F_4$ , and  $\tilde{Z}_2$  and  $\tilde{Z}_3$  are the corresponding multiplicative coefficients of the ghost-ghost wave function term  $\bar{\eta}\partial^2\eta$  and ghost-ghost-field coupling  $\partial\bar{\eta}A\eta$ . Since in the initial expression the power of the coupling constant is equal to the degree of the term minus two, relations (III.5.22) automatically express this rule and prove that up to a rescaling of the fields there is only a single coupling constant renormalization.

Relations (III.5.22) were initially proved “by hand” (see e.g. [LZ]) using so called Slavnov-Taylor identities [Sla] [Tay]. The modern method is by use of BRS symmetry [BRS], using dimensional regularization, which does not break the symmetry. A correct use of dimensional regularization remains a delicate task, however [BeDa2], and more important, it has no functional counterpart up to now. In our case, the regularization is not dimensional, and we prefer to return to the old fashioned analysis of [LZ]. In this point of view the relations (III.5.22) are deduced from some identities on Schwinger functions. At the formal level one introduces the generating functional for the theory with gauge condition  $\frac{\alpha}{2}(\partial A)^2$  (the case  $\alpha = 1$  corresponds to the Feynman gauge):

$$W(J) = \langle e^{-F^2 - (\alpha/2)(\partial A)^2} \det K e^{J \cdot A} \rangle \quad (\text{III.5.23})$$

where  $K = \partial \cdot D$  is the Faddeev-Popov operator, and the expectation value is with respect to the formal Lebesgue measure. As in [LZ][IZ], we write  $G = K^{-1}$ , where it is understood that in  $K$  and  $K^{-1}$  the  $A$  field should be replaced by the corresponding  $\frac{\delta}{\delta J}$  derivation. Then in (III.5.23) we perform a change of variables  $A \rightarrow A + D\gamma$  with  $\gamma = K^{-1}\omega$ ; by (infinitesimal) gauge invariance, there is no first order dependence in  $\omega$ , which gives the Ward-Takahashi-Slavnov-Taylor-... equation (omitting group indices):

$$\left( \alpha \frac{\partial}{\partial x^\mu} \frac{\delta}{\delta J_\mu(x)} + \int d^4y J_\lambda(y) D_\lambda(y) K_{y,x}^{-1} \right) W(J) = 0 \quad (\text{III.5.24})$$

We apply an other  $\frac{\delta}{\delta J}$ , take a divergence and put  $J = 0$  to get an identity on the two point function:

$$\langle \frac{\partial}{\partial x^\mu} A_\mu^a(x) \frac{\partial}{\partial y^\nu} A_\nu^b(y) \rangle = \frac{\delta^{ab}}{\alpha} \delta(x-y) \quad (\text{III.5.25})$$

Since in a general gauge with parameter  $\alpha$  the propagator is diagonal in su(2) space and is equal, in Fourier space, to:

$$(k^2 \delta_{\mu\nu} - (1-\alpha) k_\mu k_\nu)^{-1} = \frac{\delta_{\mu\nu} k^2 - k_\mu k_\nu}{(k^2)^2} + \frac{k_\mu k_\nu}{\alpha(k^2)^2} \quad (\text{III.5.26})$$

the first term is killed by  $-\partial_\mu\partial_\nu = k_\mu k_\nu$ , hence (III.5.25) expresses both that there is no mass renormalization and that the parameter  $\alpha$  (the gauge condition) does not renormalize.

Differentiating two more times and taking divergences before putting  $J = 0$  leads to a particularly simple identity analogous to (III.5.25) but involving the four point function, which (omitting all indices) takes the form:

$$\langle \partial A \partial A \partial A \partial A \rangle^T = 0 \quad (\text{III.5.27})$$

where the truncation corresponds to the subtraction of three terms of the type  $\langle \partial A \partial A \rangle \langle \partial A \partial A \rangle$ . It is possible to grasp intuitively that relation (III.5.27) leads to the first of the magic relations (III.5.22), namely to  $Z_2 Z_4 = Z_3^2$  [LZ]; indeed in terms of one particle irreducible components (vertices) the truncated four point function piece in (III.5.27) may be pictured as in Fig III.5.4, with the first term proportional to the four point vertex  $Z_4$  and the other ones proportional to  $Z_3^2 Z_2^{-1}$  (because they are made of two three point vertices and one propagator).

Finally one needs to derive from (III.5.24) a last relation, involving the three point function (unfortunately slightly more complicated), to obtain the second magic relation in (III.5.22), which ensures that the renormalization of the ghost-ghost-field vertex remains synchronized with the rest. It is ([LZ]):

$$\langle \left( \frac{\partial}{\partial x^\mu} A_\mu(x) \frac{\partial}{\partial y^\nu} A_\nu(y) A_\lambda(z) \right) + \frac{1}{\alpha} D_{\lambda,z} K_{zx}^{-1} \frac{\partial}{\partial y_\nu} A_\nu(y) \rangle = 0 \quad (\text{III.5.28})$$

The strategy for computing the effective expansion for the well defined theory (III.5.20-21) is to find the rigorous analogue of these identities for the theory with external (background) field  $A_i$  and internal field  $A_\rho - A_i = \sum_{j=i+1}^\rho$  integrated with the corresponding gaussian measure. This theory has an ultraviolet cutoff at scale  $\rho$  and an infrared cutoff at scale  $i$ ,  $i = \rho - 1, \rho - 2, \dots, 0$ .

In each slice there is both a gauge breaking ultraviolet and infrared cutoff. Therefore the relations (III.5.25) and (III.5.27-28) are no longer exact. Order by order in perturbation the violations can be studied in terms of Feynman diagrams (like is done at one loop in the previous section). These violations fall into two classes, the ones linked to the ultraviolet cutoff and the ones linked to the infrared cutoff. The counterterms  $CT$  in (III.5.21) have been fitted to cancel the violations of gauge invariance due to the ultraviolet cutoff at scale  $\rho$ . The violations in the  $\rho$  slice due to the infrared cutoff have an opposite sign; the corresponding difference generates the discrete flow for the counterterms  $CT$ . All violating effects which are not included in  $CT$  correspond to contributions which are convergent from the point of view of power counting. In this way apart from the flow of the two wave functions constants corresponding to the  $A$  field and ghost fields rescalings, there is a single flow for the effective coupling constant.

In the end at the zero slice we obtain a renormalized coupling constant  $\lambda_{ren}$ . Like in section II.4 we can invert the relation between  $\lambda_\rho$  and  $\lambda_{ren}$ , reexpress everything in the theory as a formal power series in  $\lambda_{ren}$  and let  $\rho$  tend to infinity. In the case of a theory on a compact torus we can then check the Slavnov Taylor identities (III.5.25) and (III.5.27-28) order by order in the renormalized coupling; in the limit  $\rho \rightarrow \infty$  all violations disappear because they correspond to operators which are irrelevant, hence exponentially small in  $\rho$ . This means that at the

perturbative level we have the right renormalized theory and that gauge invariance (in the form of Taylor-Slavnov identities) has been recovered in the limit  $\rho \rightarrow \infty$ , as expected.

Furthermore the analysis of section II.3 leads to uniform bounds of the type  $c^n \sqrt{n!}$  at order  $n$  on the renormalized series in  $\lambda_{ren}$  (the square root corresponding to the fact that the quartic couplings have couplings  $\lambda^2$  rather than  $\lambda$  like in  $\phi^4$ ).

This derivation is perhaps too sketchy to be considered a full proof, but we do not see any basic difficulty in implementing this program in more details. In this way both the standard result on renormalization of non-abelian gauge theories and a concrete bound on  $n$ -th order contributions can be derived, and what is probably more important, an effective theory with only useful renormalization performed can be derived, together with a discrete flow equation for the running coupling constant. These are however only results in the sense of formal power series, and for instance although it is clear that the first term in this flow equation is asymptotically equal to the usual  $\beta_2$  coefficient which has the correct sign for ultraviolet asymptotic freedom, it is impossible to check this asymptotic freedom in the sense of a behavior like (III.5.12) corresponding to a finite  $\lambda_{ren}$ .

### E) The positivity and domination problems

What happens if we try to apply directly the constructive method (a multi-scale expansion) to the bare theory (III.5.20-21)?

We meet problems related to a lack of positivity of the functional integrals both for a direct bound on functional integration and when trying to dominate low momentum fields produced by the expansion.

The direct lack of positivity can be studied even in a single slice model in which there is no domination problem. Unfortunately the interaction terms in (III.5.20-21) are not all of the form of positive even monomials, and this is the main difference with a  $\phi^4$  theory such as the one of section III.3. The interaction term  $F_3$  in particular is not positive, and the most natural positive bound recombines it with the gaussian piece  $F_2$  and the quartic piece  $F_4$  to reconstruct  $F^2$ . Then only the gauge fixing term does as a consequence the one loop  $A^4$  contribution behaves as

$$\left( -\frac{36 \times 2}{4} + \frac{90 \times 6}{8} - \frac{54 \times 14}{16} \right) \log K + \text{finite terms} = 2.25 \log K + \text{finite terms} \quad (\text{III.5.19})$$

where "finite terms" now means terms which are uniformly bounded both as  $\rho$  and  $K$  tends to  $+\infty$ . Therefore taking  $K$  large enough (depending on the details of our cutoff, which are responsible for the particular value of the finite terms) we can always achieve our goal of a positive total  $A^4$  contribution, hence of a negative stabilizing counterterm.

In conclusion our starting ansatz for the bare theory is: remain, and it is obviously not positive definite, hence not sufficient for functional integration. In small field regions ( $A \ll \lambda^{-1}$ ), we can bound the  $F_3$  term by the small field condition. But then we need a non perturbative bound in the regions  $A \simeq \lambda^{-1}$  which tells us that the functional weight of these regions is small compared to the (gaussian) weight of the region  $A \simeq 0$ . Although the  $A^4$  term shows that the corresponding functional weight is bounded, we cannot prove solely with this term that this weight is small.

This is already in the author's opinion a serious problem, but presumably one can, in a single slice model, find some palliative solution, for instance using the fact that a typical slice cutoff of fixed finite width  $M$  enhances the gaussian measure by a finite factor without enhancing similarly the interaction, so that it should be possible to prove that for such a fixed width cutoff the normalization of the theory tends to the gaussian normalization as  $\lambda \rightarrow 0$ . Rather than further discussing this point we prefer to focus on the domination problem, which is perhaps more interesting because it is certainly intimately related to the decomposition of the theory into several scales.

A phase space expansion consists in both spatial cluster expansions and momentum decoupling expansions. The only important restriction to build these interpolations is to preserve positivity requirements for the interpolated theory. The covariance (III.5.5) in the Feynman gauge is the same as the one of infrared  $\phi_4^4$ , hence for the horizontal expansion there is no particular positivity problem and we can use the tree cluster expansion scheme of the previous sections.

For the vertical  $t$  decoupling expansion, we need to interpolate in a way which preserves positivity and allows "domination" as much as possible. For an even positive monomial  $\lambda A^{2n}$ , let us consider the field  $A$  to be the sum  $H + L$ , where  $H$  and  $L$  represent the high and low momentum parts of  $A$ . The interpolations  $(H + tL)^{2n} + (1 - t^{2n})L^{2n}$  or  $t(H + L)^{2n} + (1 - t)L^{2n}$  are both positive and suited for domination.

As remarked above, the interaction terms in (III.5.20-21) (in particular the  $F_3$  term) are not of this form. Our first task is therefore to distinguish, in the exponential of the interaction, between two types of terms. In the first category are the ones which we call "dominable". These are the ones which come from the decomposition of even positive monomials like above, or which cannot really couple high and low momenta without violating momentum conservation, or which contain only anticommuting low momentum fields (ghosts); such fields can indeed be bounded using the Pauli principle, as in section III.4. In the second category we put the rest. Arguments based on conservation of momenta work only if there is a gap of at least one momentum slice. Therefore let us cut the propagator into the  $\rho$ -th slice, the  $\rho - 1$  slice and the rest with a  $C_0^\infty$  momentum cutoff. Explicitly we write

$$A_\rho = A^\rho + A^{\rho-1} + A_{\rho-2} \quad (\text{III.5.29})$$

$$\bar{\eta}_\rho = \bar{\eta}^\rho + \bar{\eta}^{\rho-1} + \bar{\eta}_{\rho-2} \quad (\text{III.5.30})$$

$$\eta_\rho = \eta^\rho + \eta^{\rho-1} + \eta_{\rho-2} \quad (\text{III.5.31})$$

so that (III.5.20) becomes:

$$\exp(J_{\rho-1} + I + K) d\mu_\rho(A_\rho) d\nu_\rho(\bar{\eta}_\rho, \eta_\rho) \quad (\text{III.5.32})$$

where  $I$  are the dominable terms,  $J_{\rho-1}$  is the part of the interaction which contains only low momentum fields, and  $K$  is the rest, or "non-dominable" terms. In  $I$  we want to put first all the coupling pieces coming from counterterms, because they are all dominable by the  $A^4$  term, which is a positive monomial. So we define:

$$I_1 = CT(A_\rho) - CT(A_{\rho-1}) \quad (\text{III.5.33})$$

We put also in  $I$  terms which will be treated with the Pauli principle:

$$I_2 = \lambda \{ \partial \bar{\eta}_\rho [A^\rho, \eta_\rho] + \partial \bar{\eta}^\rho [A^{\rho-1}, \eta_{\rho-1}] + \partial \bar{\eta}_{\rho-1} [A^{\rho-1}, \eta^\rho] \} \quad (\text{III.5.34})$$

Then we have terms for which gaussian integration will suffice:

$$I_3 = \lambda \{ \partial \wedge A_\rho [A^\rho, A^\rho] + \partial \wedge A_{\rho-1} ([A^\rho, A^{\rho-1}] + [A^{\rho-1}, A^\rho]) + \partial \wedge A^\rho [A^{\rho-1}, A^{\rho-1}] \} \quad (\text{III.5.35})$$

Finally we have dominable terms coming from the double commutator:

$$\begin{aligned} I_4 = \lambda^2 & \{ ([A^\rho + A^{\rho-1}, A^\rho + A^{\rho-1}] [A^\rho + A^{\rho-1}, A^\rho + A^{\rho-1}] - [A^{\rho-1}, A^{\rho-1}] [A^{\rho-1}, A^{\rho-1}]) \\ & + ([A_{\rho-2}, A^\rho + A^{\rho-1}] [A^\rho + A^{\rho-1}, A^\rho + A^{\rho-1}] - [A_{\rho-2}, A^{\rho-1}] [A^{\rho-1}, A^{\rho-1}]) \\ & \quad + \text{permutations of } A_{\rho-2} \} \\ & + 2[A_{\rho-2}, A_{\rho-2}] ([A^\rho + A^{\rho-1}, A^\rho + A^{\rho-1}] - [A^{\rho-1}, A^{\rho-1}]) \} \end{aligned} \quad (\text{III.5.36})$$

All other terms are non dominable and put into  $K$ . Basically  $K$  contains the terms with commutators between high and low momentum fields and the coupling between two high momentum ghosts and a low momentum field. More explicitly we have:

$$\begin{aligned} K = \lambda \partial \wedge A^\rho & ([A^\rho, A_{\rho-1}] + [A_{\rho-1}, A^\rho] + [A^{\rho-1}, A_{\rho-2}] + [A_{\rho-2}, A^{\rho-1}]) \\ & + \lambda \partial \wedge A^{\rho-1} ([A^\rho, A_{\rho-2}] + [A_{\rho-2}, A^\rho]) \\ & + 2\lambda^2 \left( [A^\rho + A^{\rho-1}, A_{\rho-2}] ([A^\rho + A^{\rho-1}, A_{\rho-2}] + [A_{\rho-2}, A^\rho + A^{\rho-1}]) \right. \\ & \quad \left. - [A^{\rho-1}, A_{\rho-2}] ([A^{\rho-1}, A_{\rho-2}] + [A_{\rho-2}, A^{\rho-1}]) \right) \\ & + \lambda (\partial \bar{\eta}^\rho [A_{\rho-2}, \eta^\rho + \eta^{\rho-1}] + \partial \bar{\eta}^{\rho-1} [A_{\rho-2}, \eta^\rho]) \end{aligned} \quad (\text{III.5.37})$$

If the reader adds  $I = I_1 + I_2 + I_3 + I_4$ ,  $J_{\rho-1}$  and  $K$  algebraically he will not find all of  $J_\rho$ , but all the missing terms couple one single field of slice  $\rho$  to two or three fields of cutoff  $\rho-2$ . Hence since the cutoff has compact support these terms are identically 0 when integrated over the torus  $\Lambda$  (provided  $M$  is large enough).

The non dominable terms in  $K$  create a problem. However they have quite a regular structure. If one consider the low momentum field as a background field, they come from replacing the ordinary derivatives by covariant derivatives with respect to this background field. This suggest that one should adapt the gaussian measure (and in particular the Feynman gauge condition) in the higher slices to the low momentum background field. We tried to follow this idea without much success until now. To implement this idea is particularly difficult if, like here, the gaussian measure itself is used to create the momentum slices, but even if an other method is used (such as Balaban's averages) the problem remains difficult.

This resistance of the gauge fixed non abelian functional integral to ordinary treatments is presumably related to the phenomenon of gauge ambiguity discovered by Gribov, which establishes that not every gauge orbit intersects the gauge condition once. Gribov discovery can be considered also as expressing a lack of

positivity and monotonicity of the gauge fixed functional measure. Therefore our last section is devoted to a brief discussion of this problem. Its analysis points towards abandoning the simple ansatz (III.5.20-21) for a more sophisticated one in which additional gauge conditions resolve the Gribov ambiguities. This seems to imply both surprising effects and very complicated formulas, perhaps untractable with present day constructive technology.

## F) The Gribov problem

Gribov discovered [Gri] that in the Landau gauge there can be different smooth field configurations which are nevertheless related by a gauge transformation. A configuration  $A_2 \neq A_1$  such that  $\partial_\mu A_{\mu 1} = \partial_\mu A_{\mu 2} = 0$  and such that there exists a gauge transformation  $g$  with  $A_1 = A_2^g$  ( $A^g$  being defined as in (III.5.3)) is called a Gribov copy of  $A_1$ . The weak Gribov phenomenon is that there are always some configurations which have copies; it is true even on a compact space and for configurations with the same topological properties, hence inside a given topological sector; it is also true for any regular gauge [Sin], not just the landau gauge. What we call the strong Gribov phenomenon is when there are Gribov copies of the 0 configuration, hence pure gauges which satisfy the gauge condition. This is possible in an infinite space if weak decay at infinity is allowed, as shown in [Gri], but typically this strong Gribov phenomenon does not occur inside a given topological sector in a compact space, or under strong decay conditions at infinity if the space is not compact. For instance we have:

**Lemma III.5.2** Let  $A$  be a pure gauge which is a smooth configuration vanishing outside some compact domain  $\Delta$ . Then if  $A$  satisfies the Landau condition,  $A$  is identically zero.

**Proof** We write:

$$0 = \sum_\nu \int_\Delta d^4x \int_0^{x^\nu} dy^\nu A_\nu \cdot (\sum_\mu \partial_\mu A_\mu)(x(y^\nu)) \quad (\text{III.5.38})$$

where by definition  $x^\mu(y^\nu) = x^\mu$  for  $\mu \neq \nu$ ,  $x^\nu(y^\nu) = y^\nu$  (and the dot is a scalar product in  $\text{su}(2)$ ). Integrating by parts, if  $\nu \neq \mu$  we find:

$$\begin{aligned} & \sum_\mu \sum_{\nu \neq \mu} - \int_\Delta d^4x \int_0^{x^\nu} dy^\nu \partial_\mu A_\nu \cdot A_\mu(x(y^\nu)) \\ &= - \sum_\mu \sum_{\nu \neq \mu} \int_\Delta d^4x \int_0^{x^\nu} dy^\nu (1/2) \partial_\nu (A_\mu^2)(x(y^\nu)) \\ &= - \frac{3}{2} \sum_\mu \int_\Delta d^4x (A_\mu^2)(x) \end{aligned} \quad (\text{III.5.39})$$

where in the first equality we used that  $A$  is a pure gauge, hence  $F_{\mu\nu} = 0$ , which gives  $\partial_\mu A_\nu \cdot A_\mu = \partial_\nu A_\mu \cdot A_\mu$ .

If  $\nu = \mu$  we get instead:

$$\sum_\mu \int_\Delta d^4x \int_0^{x^\mu} dy^\mu A_\mu \cdot \partial_\mu A_\mu(x(y^\mu)) = \frac{1}{2} \sum_\mu \int_\Delta d^4x A_\mu^2(x) \quad (\text{III.5.40})$$

Combining both equations we obtain  $-\int_{\Delta} \sum_{\mu} A_{\mu}^2 = 0$ , hence  $A$ , being smooth is identically zero.

The strong Gribov phenomenon happens when a change of topological sector or weak decay at infinity is allowed. In the original paper [Gri] the following example is considered. Let us distinguish the time coordinate ( $\mu = 0$ ) from the three spatial coordinates, for which we use roman letters like  $i$  or  $j$ . In the (simpler) case of a Coulomb gauge ( $A_0 = 0$  plus the condition  $\partial_j A_j = 0$ ) a pure gauge is constant in the time coordinate ( $F_{0,j} = 0 \Rightarrow \partial_0 A_j = 0$ ). The problem is then to find a three dimensional smooth configuration  $A_j$  with  $\partial_j A_j = 0$  which is a pure gauge, hence satisfies to  $A_j = \partial_j g \cdot g^{-1}$  where  $g$  is a smooth function from  $\mathbb{R}^3$  to  $SU(2)$ , the three dimensional sphere in quaternionic space. In the usual quaternionic coordinates we can write

$$g = \cos \alpha + \Omega \sin \alpha \quad (\text{III.5.41})$$

where  $\Omega$  is the vector piece of  $g$  (hence a unit three dimensional vector combination of Pauli matrices). Then

$$A_j = (\partial_j \alpha) \Omega + \sin \alpha \cos \alpha \partial_j \Omega + \sin^2 \alpha \partial_j \Omega \wedge \Omega \quad (\text{III.5.42})$$

The gauge condition in terms of  $g$  gives:

$$\begin{aligned} \partial_j A_j &= \Omega \Delta \alpha + \sin \alpha \cos \alpha \Delta \Omega + 2 \cos^2 \alpha \partial_j \alpha \partial_j \Omega \\ &= \partial_j (\Omega \partial_j \alpha + \frac{\sin 2\alpha}{2} \partial_j \Omega) \end{aligned} \quad (\text{III.5.43})$$

In the case of spherical symmetry, i.e.  $\Omega_k = x_k/r$  and  $\alpha$  solely function of  $r$ , the gauge condition becomes the equation:

$$\frac{d^2 \alpha}{dr^2} + \frac{2}{r} \frac{d\alpha}{dr} - \frac{1}{r^2} \sin 2\alpha = 0 \quad (\text{III.5.44})$$

or using  $u = 2\alpha$  and  $t = \log r$ :

$$\frac{d^2 u}{dt^2} + \frac{du}{dt} - 2 \sin u = 0, \quad (\text{III.5.45})$$

the equation of a damped pendulum. There is a solution, smooth at  $r = 0$ , which tends to  $u = \pi$  or  $\alpha = \pi/2$  for  $t \rightarrow \infty$ , hence decreases like  $1/r$ . This solution is therefore a Gribov copy of the origin (strong Gribov phenomenon).

We think that in general there is no strong Gribov phenomenon for smooth configurations on a smooth compact 4-dimensional space. Let us prove it by hand in the case of a smooth configuration in the Landau gauge on the (four dimensional) torus  $T_4$ . Let us integrate by parts the expression

$$I = \int_{T_4} (\partial_{\mu} A_{\mu}) \cdot (x^{\nu} \Omega_{\nu} \tan \alpha + x^{\nu} \alpha_{\nu} \Omega) \quad (\text{III.5.46})$$

where the product is a scalar product in quaternionic space. This expression is zero by the gauge condition. On the other hand, using (III.5.42) it is equal to

$$I = \int_{T_4} x^{\mu} \alpha_{\mu} \alpha_{\nu\nu} - x^{\mu} \left( \frac{\sin^2 \alpha}{2} \right)_{\mu} \Omega_{\nu} \Omega_{\nu}$$

$$+ \sin^2 \alpha x^\mu \Omega_\mu \Omega_{\nu\nu} + 2\alpha_\nu \cos \alpha \sin \alpha x^\mu \Omega_\mu \Omega_\nu \quad (\text{III.5.47})$$

where we adopt the notation  $f_\mu = \partial_\mu f$ . Integration by parts on a torus does not give boundary terms. Hence we get (in dimension 4):

$$I = \int_{T_4} \sum_\mu [(\alpha_\mu)^2 + \sin^2 \alpha (\Omega_\mu)^2] \quad (\text{III.5.48})$$

If  $I = 0$  we must have  $\alpha$  constant. If this constant is  $k\pi$ , then  $g = 1$  and  $A = 0$ . If it is not equal to  $k\pi$  we must have  $\Omega$  constant, hence  $g$  constant and again  $A = 0$ .

In constructive theory we are interested into explicit proofs of absence of the strong Gribov phenomenon such as this one or Lemma III.5.2 because it might point the way to useful inequalities showing that there is some definite positivity which lies in the combination of the action  $F^2$  and the Feynman gauge term  $(\partial_\mu A_\mu)^2$ . Unfortunately the positivity which comes from (III.5.39-40) or (III.5.48) seems too weak when the frequencies considered are much larger than the inverse size of the volume cutoff. In intuitive terms, the positivity which prevents the strong Gribov phenomenon is tied to boundary conditions: it is useful for the last (physical) momentum slices in a phase space analysis but does not seem strong enough at high momenta.

Although in the constructive study of the ultraviolet limit of non abelian gauge theories we can avoid the strong Gribov phenomenon just as we avoid the infrared problem (i.e. by appropriate boundary conditions or compactification), we cannot avoid the weak Gribov phenomenon. This phenomenon indeed certainly occurs in the vicinity of null vectors of the Faddeev-Popov operator

$$K = -\partial_\mu D_\mu = -\Delta + \lambda \partial_\mu [A_\mu, .] \quad (\text{III.5.49})$$

For  $\lambda = 0$  the Faddeev Popov operator reduces to minus the laplacian and is positive definite once the 0-momentum mode (hence translation invariance) has been deleted. But for non zero  $\lambda$  and  $A$  it is possible to show that rescaling  $A \rightarrow kA$  one can always have negative eigenvalues of  $K$  for  $k$  large enough (which correspond physically to bound states of the ghost system). This fact is intimately linked to the Gribov phenomenon. The configurations where  $\det K = 0$ , hence where there exists null vectors for  $K$  are the so-called Gribov horizons. The regions inside the first Gribov horizon, where  $K$  is positive is called the first Gribov region and so on. Then in [Gri] it is shown that near a Gribov horizon there are typically Gribov copies, one on each side of the horizon. These copies can be rapidly decreasing at infinity (or smooth on a compact space) so this “weak Gribov phenomenon” is completely general and cannot be eliminated like the strong one by an infrared cutoff or by topological restrictions.

Remark however that in terms of functional integration the weak Gribov phenomenon is less dangerous than the strong one. Roughly speaking we can think to the strong Gribov phenomenon as a lack of strict positivity and to the weak Gribov phenomenon as a lack of monotonicity for the functional measure. The fact that the Faddeev-Popov operator is not always positive definite at large fields is nevertheless quite disturbing. Since the determinant of this operator occurs in the functional measure, we must conclude that the ordinary measure formula for functional integration is a signed measure.

It is argued in [Hi] that this signed measure, although derived in an incorrect way, is nevertheless the correct one; essentially the argument is that as we move away from the gauge orbit of the origin (which, by absence of the strong Gribov phenomenon in our case, cuts the gauge condition only once) the Gribov copies of the weak Gribov phenomenon should by some sort of continuity argument occur in pairs with equal and opposite values, therefore in pairs which cancel out so that the ordinary prescription with the signed determinant is equivalent to a more correct prescription in which a single point on each gauge orbit is selected, and the absolute value of the determinant of the operator  $K$  is used. In any case even if this conjecture was true it seems difficult to prove it; in constructive theory one does not know well how to use signed measures and to check rigorously the corresponding cancellations, especially if they are not a marginal effect but are truly crucial for the existence of the limit.

Thinking over to the difficulties in positivity and domination that are met in the constructive analysis, and presumably their link to the existence of the weak Gribov phenomenon, it is tempting to conclude that the Gribov phenomenon should be eliminated at the beginning by use of a better ansatz than (III.5.20-21). We will limit ourselves here to a brief discussion of the corresponding approach of Zwanziger and Dell'Antonio [Zw1-3] [DeZw], in which one tries to define first a correct configuration space for the functional measure, eliminating Gribov copies by additional gauge conditions. The simplest condition which comes to mind is to choose on gauge orbits the points closest to the origin. The simplest norm to measure distance to the origin is the  $L^2$  norm on each component:

$$\|A\| = (\int Tr A^2)^{1/2} = [\int \sum_{\mu,a} (A_{\mu}^a)^2]^{1/2} \quad (\text{III.5.50})$$

In [Zw1] it is shown that on a gauge orbit the condition that  $A$  is a critical point for the function  $A \rightarrow \|A\|$  is the Landau gauge condition and the condition that it is a second order local minimum is the condition that at  $A$  the operator Faddeev-Popov operator  $K(A)$  in (III.5.49) is positive. This remarkable fact is obtained by a very simple computation. At first order, using integration by parts:

$$\int (A^{\epsilon})^2 = \int Tr A^2 - 2 \int (\partial_{\mu} A_{\mu} \cdot \epsilon) \quad (\text{III.5.51})$$

and if we use the Landau gauge condition, the second order term at a critical point is:

$$\int (\partial \epsilon \cdot \partial \epsilon + \epsilon [A, \partial \epsilon]) = - \int \epsilon K \epsilon \quad (\text{III.5.52})$$

Let us consider the configuration space  $\Lambda$  obtained by selecting the absolute minimum of the norm  $\|\cdot\|$  is selected on each gauge orbit (plus additional conditions if there are several absolute minima, which is hopefully not a generic case). We conclude that  $\Lambda$ , which should be a correct configuration space for functional integration, lies entirely inside the first Gribov region, in which the operator  $K$  is positive. In this region the ordinary Faddeev-Popov term is positive, so that we should have a better constructive starting point.

The condition that the functional integral should be restricted to a domain inside the first Gribov region has surprising consequences which contradict the standard wisdom about the high energy behavior of non-abelian gauge theory.

Indeed in the Landau gauge the operator  $K$  reduces to  $-\Delta + \lambda[A_\mu, \partial_\mu] = -\Delta + L(A)$ . Following [DeZw] we rewrite the condition that  $K$  is positive as

$$L'(A) \geq -1 \quad (\text{III.5.53})$$

where  $L'(A) = \Delta^{-1/2}L(A)\Delta^{-1/2}$ . But  $L'$  is linear in  $A$  (as is  $L$ ). It is therefore (in the case of  $SU(2)$ ) the sum of three components  $L'_j$ , each of which is antisymmetric both with respect to  $x$  space variables (because of the derivative in the Fadeev-Popov coupling) and gauge indices (because of the commutator). The tensor product of two real antisymmetric operators has a spectrum which is real and symmetric with respect to the origin, hence when such an operator is larger than  $-1$  it is also smaller than  $1$ . From this observation (and orthogonality of the eigen states for different values of  $j$ ) one concludes that when an operator like  $L'$  is bigger than  $-1$  it is also smaller than  $3$  [DeZw] (similar bounds of course also exist in the case of a group larger than  $SU(2)$ ).

The inequality  $-1 \leq L'(A) \leq 3$  on an operator which is linear in  $A$  strongly suggests that  $A$  itself should be bounded in terms of a suitable norm. Indeed this is the case and in [DeZw] such a bound is found (and called an ellipsoidal bound because the norm is not the standard one (III.5.50)). It is proved that e.g. on the torus for  $SU(2)$  one has for any  $A$  in the first Gribov region

$$N(A)^2 = \sum_k \frac{\text{Tr} A_\mu A_\mu(k)}{k^2} \leq 240 \quad (\text{III.5.54})$$

where the sum  $\sum_k$  stands for a sum over the frequencies on the lattice dual to the torus, with the zero frequency deleted.

It is reasonable to select a single point on each gauge orbit to form the correct configuration space for functional quantization. Such a point can be selected by adding further restrictions to the condition of second order minimum for the  $L^2$  norm (this looks certainly more reasonable than what is done in standard textbooks, i.e. to use simply the Landau gauge condition which selects the whole class of all stationary points in the orbit, and perhaps to hope for cancellations like in [Hi]). Therefore it seems very reasonable to consider that the functional integration for the non-abelian gauge theory can be written as a positive measure supported on a subdomain of the first Gribov region. But then one finds a contradiction with the standard perturbative renormalization group high energy behavior of the theory, at least for gauge variant quantities, because averaging the pointwise inequality (III.5.54) on configuration space with a positive measure leads to a behavior of the Fourier transform of the two point function:

$$\langle A_\mu^a(x) A_\nu^b(0) \rangle = \sum_k (\delta_{\mu\nu} - \frac{k_\mu k_\nu}{k^2}) \delta^{ab} g(k) e^{ikx} \quad (\text{III.5.55})$$

with

$$\sum_k \frac{g(k)}{k^2} \leq 80/3 \quad (\text{III.5.56})$$

where in contrast the standard renormalization group analysis [IZ] predicts:

$$g(k) \simeq_{k \rightarrow \infty} \frac{1}{k^2(\log k)^{13/22}} \quad \text{hence} \quad \sum_k \frac{g(k)}{k^2} = +\infty \quad (\text{III.5.57})$$

It is unclear whether the high energy behavior of gauge invariant quantities is also modified; this would be even more surprising because the high energy behavior of objects such as expectation values of Wilson loops computed by Monte-Carlo simulations seem in good agreement with the standard perturbative asymptotic freedom.

Much work remains of course to put the work in [DeZw] on a firm constructive basis. In particular one should start with a regularized measure with an ultraviolet cutoff which when the cutoff is removed converges towards the right measure on the right configuration space, without Gribov copies. Some first steps in this program are made in [Zw2-3].

Although it may be too early to derive conclusions, it seems that further interesting and surprising results may await us in this direction. From the physical point of view we conclude that the Gribov phenomenon appears as a non perturbative effect which seems to prevent the coexistence of too many energy scales in the non-abelian theory. On the infrared side, this was already the content of the initial paper of Gribov [Gri], in which it was argued that this effect leads to an effective infrared cutoff related to confinement. The inequalities (III.5.54-56) show that in the ultraviolet direction as well there is a truncation on frequencies which is stronger than what one expects to come from standard perturbative asymptotic freedom.

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