ALGEBRAIC RENORMALIZATION AND FEYNMAN INTEGRALS IN CONFIGURATION SPACES

ÖZGÜR CEYHAN AND MATILDE MARCOLLI

Abstract. This paper continues our previous study of Feynman integrals in configuration spaces and their algebro-geometric and motivic aspects. We consider here both massless and massive Feynman amplitudes, from the point of view of potential theory. We consider a variant of the wonderful compactification of configuration spaces that works simultaneously for all graphs with a given number of vertices and that also accounts for the external structure of Feynman graph. As in our previous work, we consider two version of the Feynman amplitude in configuration space, which we refer to as the real and complex versions. In the real version, we show that we can extend to the massive case a method of evaluating Feynman integrals, based on expansion in Gegenbauer polynomials, that we investigated previously in the massless case. In the complex setting, we show that we can use algebro-geometric methods to renormalize the Feynman amplitudes, so that the renormalized values of the Feynman integrals are given by periods of a mixed Tate motive. The regularization and renormalization procedure is based on pulling back the form to the wonderful compactification and replace it with a cohomologous one with logarithmic poles. A complex of forms with logarithmic poles, endowed with an operator of pole subtraction, determine a Rota–Baxter algebra on the wonderful compactifications. We can then apply the renormalization procedure via Birkhoff factorization, after interpreting the regularization as an algebra homomorphism from the Connes–Kreimer Hopf algebra of Feynman graphs to the Rota–Baxter algebra. We obtain in this setting a description of the renormalization group. We also extend the period interpretation to the case of Dirac fermions and gauge bosons.

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1. INTRODUCTION

In [7] and [8] we started a program aimed at studying the relation between Feynman integrals in configuration space and periods and motives of algebraic varieties. We continue here this investigation by extending the previous results in different directions.

In §2 we introduce a variant of the configuration spaces we worked with in the past. Instead of working here with a different configuration space $\text{Conf}_Γ(X)$ for each Feynman graph $Γ$, we consider a single configuration space for all graphs with a fixed number of vertices. The configuration spaces themselves, and their wonderful compactifications $\mathcal{X}_H[\mathcal{M}]$, are different from the setting we previously considered also because they incorporate some assigned data of external structure of Feynman graphs, given by imposing the condition that the configurations of points are away from a prescribed locus $H$. This type of configuration spaces, with their compactifications were previously considered by Kim and Sato, [10], and we simply we recall their main properties in our setting.

In §3 we describe the Feynman rules, the resulting Euclidean Feynman amplitudes, and their unrenormalized Feynman integrals in configuration spaces, distinguishing between the real and the complex case and between the massless and the massive case.

In §4 we consider the massless and massive propagators in configuration and momentum space and their description in terms of potential theory, using Riesz and Bessel potentials and we justify in this way more rigorously the definition of the Feynman amplitudes given in the previous section. We also describe explicitly expansions of the massive Feynman amplitudes based on expansions of modified Bessel functions of the second type.

In §5 we focus on the massive real case. We extend to the massive case the method we developed in [8] to compute massless Euclidean Feynman integrals in the real case, using expansion of the Feynman amplitude in Gegenbauer polynomials and a subdivision of the chain of integration in loci constructed using acyclic orientations of the Feynman graph. We use the expansion of the massive Feynman amplitude obtained in the previous section using expansions of the Bessel function, and we show that we can further expand the terms obtained in this way in Gegenbauer polynomials so that the same method developed in [8] applies.

In §6 we consider instead the complex case. We show that this case fits into the formalism of algebraic renormalization based on Hopf algebras of Feynman graphs and Rota–Baxter algebras of weight $-1$ as the target of the regularized Feynman amplitudes, or algebraic Feynman rules. We consider a regularization procedure for the complex Feynman amplitudes, similar to the case of [8], obtained by pulling back the smooth closed differential form given by the Feynman amplitude, with singularities along the diagonals, to a form on the wonderful compactification of the configuration space, with singularities along the boundary divisor, and by further replacing it with a cohomologous algebraic differential form with logarithmic poles along the same normal crossings divisor. We then show that one can define an associative algebra of even forms with logarithmic poles along the same normal crossings divisor. We then show that one can define an associative algebra of even forms with logarithmic poles on which there is a linear operator of “pole subtraction” that satisfies the Rota–Baxter relation of weight $-1$. The regularized amplitude defines a homomorphism of commutative algebras from the Connes–Kreimer Hopf algebra of Feynman graphs to this Rota–Baxter algebra, hence we can renormalize it using the Birkhoff factorization determined by the Hopf coproduct and antipode, on one side, and the
Rota–Baxter operator on the other. We show that this renormalization method also provides a beta function for the renormalization group.

In §7 we show that the setting of §3 and §4 can be generalized to include the cases of Dirac fermions and gauge bosons.

2. Configuration space of points

Our study of an algebro-geometric regularization procedures for position space Feynman integrals and their interpretation in terms of periods is based on wonderful compactifications and deformations of configuration space of points, where the configuration spaces considered depend on the Feynman graphs of the quantum field theory.

The geometry of the wonderful compactifications of these graph configuration spaces was examined in detail in our previous work \cite{7} and \cite{8}, see also \cite{18}, \cite{19}, and \cite{20}.

In this paper, we change this point of view slightly and we consider a fixed configuration space that corresponds to the case of the complete graph, which provides an “all in one” tool to examine Feynman integrals associated to all possible graphs. This means that, for any given Feynman amplitude, one will typically be performing more blowups than what would be strictly necessary, but with the advantage of having a single space that accommodates all the amplitudes at once. This configuration space is closely related to the original Fulton–MacPherson case \cite{14}, which also was based on the case of the complete graph, but the version we consider also contains conditions aimed at encoding the data of external structure of Feynman graphs, expressed in the form of certain “separation conditions”, as in Kim–Sato, \cite{16}. We recall here the main definitions and statements that we will need in the following. The arguments are very similar to the case previously analyzed in \cite{8} and \cite{16}, so we do not reproduce them in full.

2.1. Configuration spaces with separation conditions. We consider spaces describing configurations of \( n \) distinct points on a variety that are avoid a specified subvariety.

Let \( X \) a smooth projective variety and let \( H = \bigcup_c H_c \) be a nonsingular closed proper subvariety of \( X \) where \( H_c \) are irreducible components of \( H \) and \( c \in \mathbb{k}^+ := \{1, \ldots, k, \infty\} \).

The configuration space \( F(X \setminus H, n) \) is the complement

\[
\left( X \setminus H \right)^n \setminus \left( \bigcup_{\{i,j\} \subset \mathbb{n}} \Delta_{\{i,j\}} \right),
\]

of the diagonals

\[
\Delta_{\{i,j\}} = \{(z_i | i \in \mathbb{n}) \in X^n | z_i = z_j\},
\]

associated to the pairs \( \{i,j\} \) in \( \mathbb{n} := \{1, \ldots, n\} \) in the cartesian product \( (X \setminus H)^n \).

2.2. Wonderful compactification of configuration spaces. We now describe the various steps involved in the construction of wonderful compactifications for the configuration spaces \( F(X \setminus H, n) \).

2.2.1. Diagonals and subvarieties. We start by considering the following collection of subvarieties of \( X^n \):

- Given a non-empty subset \( S \) of \( \mathbb{n} \) and a subset \( \alpha \) of \( \mathbb{k}^+ \), we have a subvariety

\[
\Delta_{\alpha,S} = \{(z_i | i \in \mathbb{n}) \in X^n | \forall i \in S, \exists c \in \alpha \text{ with } z_i \in H_c\}
\]

- For subsets \( I \subset \mathbb{n} \) with \( |I| > 1 \), we have the corresponding diagonals

\[
\Delta_I = \{(z_i | i \in \mathbb{n}) | z_i = z_j \text{ for } i,j \in I\}.
\]
2.2.2. First step: arrangements and configurations separated from $\mathcal{H}$. The subvarieties $\Delta_{\{c\}, S}$ for $c \in \mathbb{k}^+$ are isomorphic to 
$$\Delta_{\{c\}, S} \cong \mathcal{H}_c^S \times \mathcal{X}^{n \setminus S}.$$ 
The collection $\mathcal{S}_n = \{\cap_{\alpha} \Delta_{\alpha, S}\}$ of all possible intersections of subvarieties $\Delta_{\alpha, S}$ forms a simple arrangement in $\mathcal{X}^n$.

**Lemma 2.1** (Kim & Sato [16], §2.2). The set $\mathcal{G}_n = \{\Delta_{\{c\}, S} \mid c \in \mathbb{k}^+ \& S \subset \mathcal{n}\}$ is a building set for the arrangement $\mathcal{S}_n$.

2.2.3. Blow-up construction of the first step in the compactification. We introduce a partial ordering $\preceq$ on $\mathcal{G}_n$ according to the dimension of subspaces i.e., $\Delta_{\{c\}, S} \preceq \Delta_{\{c'\}, S'}$ when $\dim(\Delta_{\{c\}, S}) \leq \dim(\Delta_{\{c'\}, S'})$.

Let $\mathcal{Y}_0 = \mathcal{X}^n$ and let $\mathcal{Y}_k$ be the blow up of $\mathcal{Y}_{k-1}$ along the iterated dominant transform of those subvarieties $\Delta_{\{c\}, S}$ whose dimension is $k$. We simply set $\mathcal{Y}_k = \mathcal{Y}_{k+1}$ if there are no subspaces $\Delta_{\{c\}, S}$ of dimension $k$. The resulting sequence of blow ups
$$\mathcal{Y}_m \longrightarrow \cdots \longrightarrow \mathcal{Y}_1 \longrightarrow \mathcal{Y}_0 = \mathcal{X}^n$$
terminates when it reaches the maximal dimension of the subvarieties $\Delta_{\{c\}, S}$, that is,
$$m = \dim(\mathcal{X}) \cdot (n - 1) + \max_c(\dim(\mathcal{H}_c)).$$

Each step of this iteration procedure is independent on the order in which the blowups are performed along the (iterated) dominant transforms of the subspaces $\Delta_{\{c\}, S}$. Thus, the intermediate varieties $\mathcal{Y}_k$ in the sequence (2.5) are all well defined.

The variety $\mathcal{X}_H^\mathcal{H}[n] := \mathcal{Y}_m$ obtained through this sequence of iterated blowups is the wonderful compactification of the arrangement $\mathcal{S}_n$. It is a compactification of the configuration space of $n$ labelled (not necessarily distinct) points in $\mathcal{X}$ that keeps the points away from the subvariety $\mathcal{H}$.

2.2.4. Second step: Configurations of distinct points. Let $\tilde{\Delta}_I \subset \mathcal{X}_H^\mathcal{H}[n]$ denote the proper transform of $\Delta_I$.

**Lemma 2.2** (Kim & Sato [16], Proposition 4). The collection of all subspaces $\{\tilde{\Delta}_I \mid I \subset \mathcal{n}, |I| > 1\}$ is a building set of a subspace arrangement in $\mathcal{X}_H^\mathcal{H}[n]$.

Let $\mathcal{M}_0 = \mathcal{X}_H^\mathcal{H}[n]$ and let $\mathcal{M}_{k+1}$ be the blow up of $\mathcal{M}_k$ along the iterated dominant transform of $\tilde{\Delta}_I$ for $|I| - k$. The resulting sequence of blow ups
$$\mathcal{M}_{n-1} \longrightarrow \cdots \longrightarrow \mathcal{M}_1 \longrightarrow \mathcal{M}_0 = \mathcal{X}_H^\mathcal{H}[n]$$
does not depend on the order in which the blowups are performed along the (iterated) dominant transforms of the subspaces $\tilde{\Delta}_I$, hence the intermediate varieties $\mathcal{M}_k$ in the sequence (2.6) are all well defined. The variety $\mathcal{X}_H^\mathcal{H}[n] := \mathcal{M}_{n-1}$ is a compactification of the configuration space of $n$ labelled points in $\mathcal{X}$ that keeps the points pairwise disjoint and away from the subvariety $\mathcal{H}$. Moreover, we have the following result.

**Theorem 2.3** (Kim & Sato [16], Theorem 2). The wonderful compactifications obtained as above have the following properties:

1. The variety $\mathcal{X}_H^\mathcal{H}[n]$ is nonsingular.
(2) The boundary
\[
\mathcal{X}_\mathcal{H}[n] \setminus \left( (\mathcal{X} \setminus \mathcal{H})^n \setminus \left( \bigcup_{\{i,j\} \subset \mathbb{Z}} \Delta_{\{i,j\}} \right) \right)
\]
is the union of exceptional divisors \( D_{c,S} \) and \( D_I \), corresponding to the subvarieties \( \Delta_{c,S} \) for \( c \in \mathbb{R}^+ \), \( |S| > 0 \), and the diagonals \( \Delta_I \) for \( |I| > 1 \). Any set of these divisors intersects transversally.

(3) The intersection of boundary divisors \( D_{c_1,S_1}, \ldots, D_{c_a,S_a}, D_{I_1}, \ldots, D_{I_b} \) is nonempty if and only if they are nested.

3. Feynman integrals of Euclidean scalar QFTs

In this section, we describe Feynman amplitudes and unrenormalized Feynman integral on the configuration space \( \mathcal{F}(\mathcal{X} \setminus \mathcal{H}, n) \).

3.1. Feynman graphs. A graph \( \Gamma \) is a one-dimensional finite CW-complex. We denote the set of vertices of \( \Gamma \) by \( V_\Gamma \), the set of edges by \( E_\Gamma \), and the boundary map by \( \partial_\Gamma : E_\Gamma \to (V_\Gamma)^2 \). We write \( \partial_\Gamma(e) = \{s(e), t(e)\} \), the source and target vertices of the oriented edge \( e \). We denote the valency of a vertex \( v \) by \( |v| \).

The set of external vertices \( V_\Gamma^e \) is \( \{v \in V_\Gamma \mid |v| = 1\} \), and the set of internal vertices \( V_\Gamma^i \) is its complement \( V_\Gamma \setminus V_\Gamma^e \). The set of edges \( E_\Gamma \) is the union of the set \( E_\Gamma^e \) of external edges, that are the edges having external vertices as boundary vertices, and \( E_\Gamma^i := E_\Gamma \setminus E_\Gamma^e \) is the set of internal edges.

A looping edge is an edge for which the two boundary vertices coincide. Since the position space Feynman amplitudes would diverge everywhere in such cases, we exclude them and assume that our graphs have no looping edges.

3.2. Feynman rules. As we discussed already in \([8]\), we consider two different approaches to defining the Feynman integrals in configuration space, which we refer to here as the real case and the complex case. In the real case, one can choose the manifold \( \mathcal{X} \) to be \( \mathbb{P}^D(\mathbb{C}) \) and the physical \( D \)-dimensional spacetime manifold will be the real locus \( \mathbb{A}^D(\mathbb{R}) \) inside an affine space \( \mathbb{A}^D(\mathbb{C}) \subset \mathbb{P}^D(\mathbb{C}) \).

In the complex case spacetime is complexified. The manifold \( \mathcal{X} \) consists of a product \( \mathbb{P}^D(\mathbb{C}) \times \mathbb{A}^D(\mathbb{C}) \) and the spacetime is the middle dimensional subspace given by an affine space \( \mathbb{A}^D(\mathbb{C}) \times \mathbb{A}^0 \) lying in one of the two copies of \( \mathbb{P}^D(\mathbb{C}) \). The two choices lead to completely different methods of dealing with the resulting Feynman integrals, due to the fact that, in the real case, the differential form describing the Feynman amplitude is not a closed form, while in the complex case it is. We recall below the two cases and the resulting Feynman amplitudes.

3.2.1. Feynman rules: real case. Let \( \mathcal{X} = \mathbb{P}^D(\mathbb{C}) \) and consider the real locus \( \mathbb{A}^D(\mathbb{R}) \) in the affine part \( \mathbb{A}^D(\mathbb{C}) \) as the physical space-time manifold. We also denote the hyperplane at infinity \( \mathbb{P}^D \setminus \mathbb{A}^D \) by \( \mathcal{H}_\infty \).

Given a Feynman graph \( \Gamma \) with \( n \) internal and \( k \) external vertices and no looping edges, the Feynman rules assigns to \( \Gamma \) a Feynman amplitude as follows:

- The external vertices are labelled by fixed pairwise disjoint points \( \mathcal{H}_c = \{z_c\} \subset \mathcal{X} \) with \( c \in \mathbb{R} \), which we refer to as “external points”;
- To each internal vertex \( v_j \) one assigns a variable \( z_j \in \mathcal{X} \), so that internal vertices are labelled by the set of points \( \{z_1, \ldots, z_n\} \in \mathcal{X}^n \);
- To each edge \( e \in E_\Gamma \) with \( \partial_\Gamma(e) = \{s(e), t(e)\} \), we assign a Euclidean propagator. This is given in the massless case by

\[
G_{0,R}(x_{s(e)} - x_{t(e)}) = \frac{1}{\|x_{s(e)} - x_{t(e)}\|^{2\lambda}}, \quad \text{where} \quad D = 2\lambda + 2.
\]
In the massive case the propagator has the form
\begin{equation}
G_{m,R}(x_{s(e)} - x_{t(e)}) = (2\pi)^{-(\lambda+1)} m^{\lambda} \| x_{s(e)} - x_{t(e)} \|^{-\lambda} K_\lambda(m \| x_{s(e)} - x_{t(e)} \|),
\end{equation}
where $K_\lambda(z)$ is the modified Bessel function (Macdonald function, see §1.2 of [1]).

We will explain more in detail the origin of the propagators (3.1) and (3.2) in terms of potential theory in [4] below, using the fact that the propagator (3.1) is obtained by considering the Green function of the real Laplacian on $A^D(\mathbb{R})$, which for $D = 2\lambda + 2$ is given by
\begin{equation}
G_{\mathbb{R},\Delta}(x-y) = \frac{1}{\| x-y \|^{2\lambda}}.
\end{equation}

The resulting Feynman amplitude associated to the graph $\Gamma$ in the massless real case is then the differential form
\begin{equation}
\omega_{\Gamma,R} = \prod_{e \in E_\Gamma} \frac{1}{\| x_{s(e)} - x_{t(e)} \|^{2\lambda}} \wedge dx_v.
\end{equation}
This is a smooth differential form in degree equal to the middle dimension, on the configuration space, $\omega_{\Gamma,R} \in \Omega^{Dn}(\mathcal{F}(\mathcal{X} \setminus \mathcal{H}, n))$, where we set
\begin{equation}
\mathcal{H} := \bigcup_{c \in \mathbb{K}^+} \mathcal{H}_c.
\end{equation}
However, it is not a closed form, see Proposition 3.3 of [8]. Similarly, the Feynman amplitude associated to the graph $\Gamma$ in the massive real case is
\begin{equation}
\omega_{\Gamma,m,R} = (2\pi)^{-(\lambda+1)|E_\Gamma|} \prod_{e \in E_\Gamma} m^{\lambda} \| x_{s(e)} - x_{t(e)} \|^{-\lambda} K_\lambda(m_e \| x_{s(e)} - x_{t(e)} \|) \wedge dx_v,
\end{equation}
where we allow the different edges to have different mass parameters, with $m = (m_e)_{e \in E_\Gamma}$.

3.2.2. Chain of integration: real case. In this real case, the chain of integration for the Feynman integral is given by the middle dimensional cycle
\begin{equation}
\sigma_{\Gamma,R} = \mathcal{X}(\mathbb{R})^V = \mathcal{X}(\mathbb{R})^n.
\end{equation}

3.2.3. Feynman rules: complex case. Let $\mathcal{X} = \mathbb{P}^D(\mathbb{C}) \times \mathbb{P}^D(\mathbb{C})$ with projection $p : \mathcal{X} \to \mathbb{P}^D(\mathbb{C})$, $p : z = (x,y) \mapsto x$. The complex affine space $A^D(\mathbb{C})$ in $p(\mathcal{X})$ is identified as the physical space-time in this case.

- The external vertices are labelled by fixed pairwise disjoint points $\mathcal{H}_c = \{ z_c = (x_c,y) \} \subset \mathcal{X}$, with $c \in \mathbb{K}$, and with $y$ a chosen base point in $\mathbb{P}^D(\mathbb{C})$.
- To each internal vertex $v_j$ one assigns a variable $x_j = p(z_j)$, with $z_j \in \mathcal{X}$, so that internal vertices are labelled by the set of points $(x_1, \ldots, x_n) \in \mathbb{P}^D(\mathbb{C})^n$;
- To each edge $e \in E_\Gamma$, with $\partial(e) = \{ s(e), t(e) \}$, we consider the Euclidean massless edge propagator
\begin{equation}
G_{0,C}(x_{s(e)} - x_{t(e)}) = \frac{-(D-2)!}{(2\pi i)^D} \frac{1}{\| x_{s(e)} - x_{t(e)} \|^{2D-2}},
\end{equation}
where $\| x_{s(e)} - x_{t(e)} \| = \| p(z)_{s(e)} - p(z)_{t(e)} \|$.

The resulting massless Feynman amplitude for a graph $\Gamma$ is then of the form
\begin{equation}
\omega_{\Gamma,C} = \left( \frac{-(D-2)!}{(2\pi i)^D} \right)^{|E_\Gamma|} \prod_{e \in E_\Gamma} \frac{1}{\| x_{s(e)} - x_{t(e)} \|^{2D-2}} \wedge dx_v \wedge d\bar{x}_v.
\end{equation}
This is a closed differential form, of degree equal to the middle dimension, in the configuration space \( \omega_{\Gamma, C} \in \Omega^{2Dn}(F(X \setminus \mathcal{H}, n)), \) with \( \mathcal{H} = \cup_{c \in \mathbb{R}} \mathcal{H}_c. \)

In the massive case, with the complexified spacetime \( \mathbb{A}^D(\mathbb{C}) \), the edge propagator (3.7) is replaced by the massive edge propagator

\[
G_{m, C}(x_{s(e)} - x_{t(e)}) = (2\pi)^{-D} m^{D-1} \|x_{s(e)} - x_{t(e)}\|^{-(D-1)} K_{D-1}(m \|x_{s(e)} - x_{t(e)}\|),
\]

with \( K_{\nu}(z) \) the modified Bessel function.

Thus, in the complexified case, the massive Feynman amplitude is given by

\[
\omega_{\Gamma, m, C} = (2\pi)^{-D|E^i|} \prod_{e \in E^i} m^{D-1} \|x_{s(e)} - x_{t(e)}\|^{-(D-1)} K_{D-1}(m \|x_{s(e)} - x_{t(e)}\|) \bigwedge_{v \in V^i} dx_v \wedge d\bar{x}_v.
\]

The choice of propagators in this case arises from the fact that on \( \mathbb{A}^D(\mathbb{C}) \) the complex Laplacian

\[
\Delta = \sum_{k=1}^{D} \frac{\partial^2}{\partial x_k \partial \bar{x}_k}
\]

has Green form

\[
G_{C, \Delta}(x - y) = \frac{-(D - 2)!}{(2\pi i)^D \|x - y\|^{2D - 2}}.
\]

We will return to this more in detail in \( \S 4 \) below.

3.2.4. Chain of integration: complex case. In the complex case, the chain of integration for the Feynman integral is given by the middle dimensional cycle

\[
\sigma_{\Gamma, C} = \mathbb{P}^D(\mathbb{C})^{V^i} \times \{y\} = \mathbb{P}^D(\mathbb{C})^n \times \{y\},
\]

for a choice of a point \( y = (y_v | v \in V^i_\Gamma). \)

3.2.5. Unrenormalized Feynman integral. The unrenormalized massless Feynman integral for a Feynman graph \( \Gamma \) is then given by

\[
\int_{\sigma_{\Gamma, R}} \omega_{\Gamma, R} \quad \text{in the real case}
\]

\[
\int_{\sigma_{\Gamma, C}} \omega_{\Gamma, C} \quad \text{in the complex case},
\]

and similarly for the massive Feynman integrals.

4. Potential theory and propagators in configuration and momentum space

In this section we justify the form (3.1) and (3.2) of the propagators in terms of potential theory, and similarly for their complex counterparts. In particular, we show how the expressions (3.1) and (3.2) arise naturally from the theory of Riesz and Bessel potentials, which gives a rigorous formulation of the Fourier transform relation between momentum and configuration space propagators.
4.1. **Massless propagator and the Riesz potential.** We first recall the relation between the massless Euclidean configuration space propagator

\[ G_{0,R,D}(x_s - x_t) = ||x_s - x_t||^{2-D} \]

for \( D > 2 \) and the momentum space propagator, through the relation of (4.1) to the Riesz potential.

**Lemma 4.1.** The propagator (4.1) satisfies the Fourier transform formula

\[ \hat{(G_{0,R,D} \star \varphi)}(k) = \frac{4\pi^{D/2}}{\Gamma(\lambda)} \frac{1}{||k||^2} \hat{\varphi}(k), \]

for any test function \( \varphi \in \mathcal{S}(\mathbb{R}^D) \).

**Proof.** The function (4.1) is the Green function of the real Laplacian on \( \mathbb{R}^D \). It acts on test functions by convolution product

\[ G_{0,R,D} \star \varphi(x) = \int_{\mathbb{R}^D} \frac{\varphi(y)}{||x - y||^{D-2}} dy. \]

This is a special case of the Riesz potential (see §1.2 of [1]) \( \mathcal{I}_\alpha = C_\alpha ||x||^{\alpha-D} \) with

\[ C_\alpha = \frac{\Gamma((D-\alpha)/2)}{2^{\alpha} \pi^{D/2} \Gamma(\alpha/2)}, \]

in the case where \( \alpha = 2 \),

\[ G_{0,R,D}(x - y) = \frac{4\pi^{D/2}}{\Gamma(\lambda)} \mathcal{I}_2(x - y) \]

where \( D = 2\lambda + 2 \). The Riesz potential satisfies the Fourier transform property

\[ \hat{(\mathcal{I}_\alpha \star \varphi)}(k) = \frac{1}{||k||^\alpha} \hat{\varphi}(k), \]

hence, \( G_{0,R,D}(x - y) \) satisfies the Fourier transform formula (4.2).

Notice that \( ||k||^{2-D} \) is the massless momentum space propagator associated to a single edge carrying a momentum variable \( k \in \mathbb{R}^D \). Thus configuration space propagator (4.1) can be identified with the Fourier transform of the momentum space propagator associated to the same edge of the Feynman graph as

\[ \hat{G}_{0,R,D}(k) = \frac{4\pi^{D/2}}{\Gamma(\lambda)} \frac{1}{||k||^2}. \]

4.1.1. **Complex case.** In the complex case, the propagator (3.7) can be similarly related to the momentum space propagator, by interpreting it as the Green function (3.11) of the complex Laplacian on \( \mathbb{A}^D(\mathbb{C}) \) as \( G_{C,\Delta}(x, y) = (-D-2)!/(2\pi i)^D)G_{0,2D} \) and using again (4.3) to relate it to the massless momentum space propagator, so that we obtain

\[ G_{0,C,D}(x_{s(e)} - x_{t(e)}) = \frac{-(D-2)!}{(2\pi i)^D} \frac{4\pi^D}{\Gamma(D-1)} \mathcal{I}_2(x_{s(e)} - x_{t(e)}) = (2i)^{2-D} \mathcal{I}_2(x_{s(e)} - x_{t(e)}), \]

which satisfies

\[ \hat{G}_{0,C,D}(k) = (2i)^{2-D} \frac{1}{||k||^2}. \]
4.2. Massive propagator and Helmholtz equation. The case of the massive propagator can be similarly related via Fourier transform to the Bessel potential. Notice that we use here a different sign convention for the Laplacian with respect to Eq. [1]. Let $G_{m,R,D}(x - y)$ denote the massive propagator of Eq. [3.2] in dimension $D = 2\lambda + 2$.

**Lemma 4.2.** The massive Euclidean configuration space propagator $G_{m,R,D}(x - y)$ satisfies the Fourier transform formula

$$
\hat{(G_{m,R,D} \ast \varphi)}(k) = \frac{1}{(m^2 + \|k\|^2)} \hat{\varphi}(k)
$$

for any test function $\varphi \in \mathcal{S}(\mathbb{R}^D)$. In particular, $G_{m,R,D}$ is the Green function of the Helmholtz equation with wavenumber $m^2$.

**Proof.** The fractional integral operator $(m^2 \cdot I + \Delta)^{-\alpha/2}$, where $m \neq 0$ is the mass parameter, satisfies the Fourier transform property

$$
((m^2 \cdot I + \Delta)^{-\alpha/2} \varphi)(k) = \frac{1}{(m^2 + \|k\|^2)^{\alpha/2}} \hat{\varphi}(k),
$$

hence it can be represented by the convolution operator

$$
(m^2 \cdot I + \Delta)^{-\alpha/2} \varphi = G^{(\alpha)}_{m,R,D} \ast \varphi, \quad \text{with} \quad G^{(\alpha)}_{m,R,D}(k) = \frac{1}{(m^2 + \|k\|^2)^{\alpha/2}}.
$$

In the case $m = 1$, this is the usual Bessel potential $G^{(\alpha)}_{\alpha,R,D}$ (see §1.2 of Eq. [1]) and, in particular, when $\alpha = 2$, the Green function $G^{(2)}_{m,R,D}$ is the fundamental solution of the Helmholtz equation with wavenumber $m^2$, namely it satisfies

$$
\Delta G^{(2)}_{m,R,D} + m^2 G^{(2)}_{m,R,D} = \delta.
$$

In the case $\alpha = 2$, the Green function $G^{(2)}_{m,R,D}$ gives the massive Euclidean configuration space propagator $G_{m,D}$, since the momentum space Euclidean massive propagator is of the form

$$
\frac{1}{m^2 + \|k\|^2},
$$

with $k \in \mathbb{R}^D$ the momentum variable associated to the edge.  

4.2.1. Complex case. In the complex case, similarly, the massive propagator is a solution of the Helmholtz equation with the complex Laplacian and wavenumber $m^2$,

$$
\Delta_c G_{m,C,D} + m^2 G_{m,C,D} = \delta.
$$

It will be clear from §4.3 below that this gives the form of the propagator given in Eq. [3.9].

4.3. Massive propagator and the Bessel potential. The massive propagator $G_{m,D}$ can then be described explicitly in terms of Bessel functions as follows.

**Lemma 4.3.** The massive Euclidean propagator $G_{m,R,D}$ on $\mathbb{R}^D$ has the following integral representation

$$
G_{m,R,D}(x) = \frac{1}{(4\pi)^{\frac{D}{2}}} \int_0^{\infty} dt \ t^{-\frac{D}{2}} e^{-tm^2 - \|x\|^2/4t}.
$$
η, α are analytic in \( \hat{D} \) which gives (4.10). Since both we can give its Fourier transform (4.6) as follows

\[
G_{m,D,\alpha}(x) = \frac{1}{(2\pi)^D} \int_{\mathbb{R}^D} dk \ e^{ik \cdot x} \int_0^\infty dt \ e^{-t(\|k\|^2 + m^2)}
\]

Proof. By using the following integral form of the momentum space propagator

\[
\frac{1}{\|k\|^2 + m^2} = \int_0^\infty dt \ e^{-t(\|k\|^2 + m^2)}
\]

we can give its Fourier transform (4.6) as follows

\[
G_{m,D,\alpha}(x) = \frac{1}{(2\pi)^D} \int_{\mathbb{R}^D} dk \ e^{ik \cdot x} \int_0^\infty dt \ e^{-t(\|k\|^2 + m^2)}
\]

\[
= \frac{1}{(2\pi)^D} \int_0^\infty dt \ e^{-tm^2} \int_{\mathbb{R}^D} dk \ e^{-t\|k\|^2} \Gamma\left(\frac{D}{2}\right) \cdot e^{\frac{-\|q\|^2}{4t}}
\]

\[
= \frac{1}{(2\pi)^D} \int_0^\infty dt \ e^{-tm^2} \|q\|^2 \int_{\mathbb{R}^D} dq \ e^{-t\|q\|^2} \text{ where } q = k - \frac{ix}{2t}
\]

and obtain (4.9) by substituting \( \frac{2\pi}{2t} \frac{D}{}\frac{D}{t} = \int_{\mathbb{R}^D} dq \ e^{-t\|q\|^2} \).

In particular, this implies the following behavior of divergences.

**Corollary 4.4.** For \( x = 0 \), we have

\[
G_{m,D,\alpha}(0) = \frac{1}{(4\pi)^\frac{D}{2}} \int_0^\infty dt \ t^{-\frac{D}{2}} \ e^{-tm^2} = \frac{1}{(4\pi)^\frac{D}{2}} m^{D-2} \Gamma(1 - \frac{D}{2}),
\]

that diverges when \( D = 2\lambda + 2 \) for \( \lambda \in \mathbb{N}^+ \).

It is convenient, in order to write \( G_{m,D,\alpha} \) in terms of Bessel functions, to first encode the mass dependence in a scaling relation in terms of the propagator at \( m = 1 \).

**Lemma 4.5.** The dependence of the Green function \( G_{m,D,\alpha} \) of the operator \( G_{m,D,\alpha} = (mI + \Delta)^{-\alpha/2} \) on the mass parameter \( m \) is given by

\[
(4.10) \quad G_{m,D,\alpha}(x) = m^{D-\alpha} G_{D,\alpha}(mx),
\]

where \( G_{D,\alpha}(x) = G_{m=1}^{(\alpha)} \) is the Bessel potential on \( \mathbb{R}^D \).

Proof. This follows from the integral representation (4.9) for \( G_{m,D,\alpha}(x) \) and the integral representation

\[
G_{1,D}^{(\alpha)}(x) = a_\alpha \int_0^\infty t^{-\frac{D}{2}} \ e^{-\frac{\|x\|^2}{4t} - \frac{t}{4\pi}} dt, \text{ where } a_\alpha = \frac{1}{(4\pi)^{\alpha/2} \Gamma(\alpha/2)}
\]

for the Bessel potential (see §1.2.4 of [I]). Using the change of variables \( t = 4\pi m^2 \cdot t' \) in (4.11) gives (4.10).

The scaling (4.10) can also be seen in the Fourier representation

\[
G_{D,\alpha}(x) = \frac{1}{(2\pi)^D} \int_{\mathbb{R}^D} \frac{e^{ix\cdot\xi}}{(1 + \|\xi\|^2)^{\alpha/2}} d\xi,
\]

where, after changing variables to \( k = \xi/m \), one gets

\[
G_{m,D,\alpha}(x) = \frac{1}{(2\pi)^D} \int_{\mathbb{R}^D} \frac{m^D e^{im\cdot\xi}}{m^\alpha(1 + \|\eta\|^2)^{\alpha/2}} d\eta,
\]

which gives (4.10). Since both \( \widehat{G_{D,\alpha}} \) and \( \widehat{G_{m,D,\alpha}} \) are entire analytic functions of \( \alpha \) for each fixed \( \xi \) or \( \eta \), the corresponding distributions \( G_\alpha \) and \( G_{m,\alpha} \) have the property that \( \langle G_{D,\alpha}, \varphi \rangle \) and \( \langle G_{m,D,\alpha}, \varphi \rangle \) are analytic in \( \alpha \) for all test functions \( \varphi \in \mathcal{S} \), and the identity (4.10) extends for all \( \alpha \).

We then obtain the following expression for the massive propagator in terms of Bessel functions.
Lemma 4.6. The massive Euclidean propagator $G_{m,D}(x-y) = (2\pi)^{-\frac{D}{2}} m^{D-2} (m \|x-y\|)^{-\frac{D-2}{2}} K_{D-2}(m \|x-y\|)$, where $K_\nu(z)$ is the modified Bessel function (Macdonald function).

Proof. This can be seen from the integral representation of Lemma 4.3 or equivalently from Lemma 4.6. We use again the fact that the Bessel potential satisfies \[ (4.11) \] which implies (2.3) p. 201 and §1.2 of [1]) that it can be expressed in terms of Bessel functions as
\[ (4.13) \]

$$G^{(\alpha)}_{1,D}(x) = c_\alpha \|x\|^{-\frac{D-\alpha}{2}} K_{\frac{D-\alpha}{2}}(\|x\|),$$

with $c_\alpha = 2^{(\alpha-2)/2} (2\pi)^{D/2} \Gamma(\alpha/2)$, and with $K_\nu(z)$ the modified Bessel function.

The statement follows by applying again the change of variables $t = 4\pi m^2 \cdot t^\prime$ in the integral (4.11) and using the identities (4.13) and (4.9) for case of $\alpha = 2$. Equivalently, (4.12) follows from the scaling property of Lemma 4.5 and (4.13).

The integral representation of Lemma 4.3 also gives an estimate on the behavior of the massive propagator at infinity.

Corollary 4.7. The massive propagator $G_{m,R,D}(x)$ is smooth on $\mathbb{R}^D \setminus \{0\}$. Moreover, $G_{m,R,D}(x) \to 0$ as $\|x\| \to \infty$.

Proof. The smoothness of $G_{m,D}(x)$ is a direct consequence of the fact that the modified Bessel function $K_\nu(z)$ is analytic on $\mathbb{C}^*$ and has an essential singularity at $z = \infty$.

The Cauchy–Schwarz inequality takes the form
\[ (4.12) \]

$$|G_{m,R,D}(x)|^2 \leq \frac{1}{(4\pi)^D} \left( \int_0^\infty dt e^{-2tm^2} \right) \cdot \left( \int_0^\infty dt \ t^{-D} e^{-\frac{\|x\|^2}{2t}} \right) = \frac{1}{(4\pi)^D} \frac{1}{2m^2} \int_0^\infty dt \ t^{-D} e^{-\frac{\|x\|^2}{2t}}$$

$$= \frac{2^{D-1}}{(4\pi)^D} \frac{1}{2m^2} \|x\|^{-2D+2} \int_0^\infty dt' \ e^{-t't} t'^{D-1} = \frac{2^{D-1}}{(4\pi)^D} \frac{1}{2m^2} \Gamma(D-1) \|x\|^{-2D+2}$$

$$= \frac{2^{D-1}}{(4\pi)^D} \frac{(D-2)!}{2m^2} \|x\|^{-2D+2}$$

after the coordinate transformation $t' = \frac{\|x\|^2}{2t}$, and it simply implies that $G_{m,R,D}(x)$ in (4.9) tends to zero at infinity for $D > 1$ due to the term $\|x\|^{-2D+2}$.

This can also be seen using the integral representation (4.11) for $G^{(\alpha)}_{1,D}$ for $\alpha = 2$ and the integral representation of the Riesz potential (see §1.2.4 of [1])
\[ (4.13) \]

$$I_\alpha(x) = a_\alpha \int_0^\infty t^{\alpha-D} \ t^{\frac{D-1}{2}} e^{-\frac{\|x\|^2}{t}} \ dt,$$

which then gives an estimate (for all $0 < \alpha < D$)
\[ (4.14) \]

$$0 \leq G^{(\alpha)}_{1,D}(x) \leq I_\alpha(x).$$

4.3.1. Complex case. The same arguments used in Lemma 4.6 and Corollary 4.7 give an expression of the massive propagator in the complex case in terms of modified Bessel functions.

Lemma 4.8. The massive Euclidean propagator $G_{m,C,D}(x-y) = (2\pi)^{-D} m^{D-1} \|x-y\|^{-(D-1)} K_{D-1}(m \|x-y\|)$

In particular $G_{m,C,D}$ is smooth away from the origin, with $G_{m,C,D}(x) \to 0$ as $\|x\| \to \infty$. 
4.4. Expansions of massive Feynman amplitudes. Using the series expansions of the Bessel functions, we obtain the following formal expansion for the massive Euclidean propagator in terms of the distance function.

**Lemma 4.9.** Let $D = 2\lambda + 2$, then the massive Euclidean propagator has an expansion of the form

\[
G_{m,\mathbb{R},D}(x - y) = \frac{1}{(2\pi)}\sum_{\ell=0}^{(\lambda + 1)} \frac{(-m^2)^\ell}{\ell! \|x - y\|^{2(\lambda - \ell)}} \cdot \frac{2^\lambda \cdot (\lambda - 1)!}{2^\ell \cdot (\lambda + \ell - 1)!} \cdot (\lambda - \ell - 1)! \cdot (\lambda - 1)! \\
+ \frac{1}{(2\pi)}\sum_{\ell=1}^{\infty} \frac{m^{2(\lambda + \ell)} \|x - y\|^{2\ell}}{2^{2\ell} \cdot \ell! \cdot (\lambda + \ell)!} \cdot \left(\log\left(\frac{m \|x - y\|}{2}\right) - \frac{1}{2} \psi(\ell + 1) + \frac{1}{2} \psi(\lambda + \ell + 1)\right)
\]

where $\psi(x) = \frac{d}{dx} \ln \Gamma(x) = \frac{\Gamma'(x)}{\Gamma(x)}$ is the digamma function.

**Proof.** We use the following expansion of modified Bessel function $K_\lambda(z)$, see p. 80 in [25],

\[
K_\lambda(z) = \frac{1}{2} \sum_{\ell=0}^{\lambda - 1} \frac{(-1)^\ell (\lambda - \ell)!}{\ell! (\frac{z}{2})^{\lambda - 2\ell}} + (-1)^{\lambda + 1} \sum_{\ell=0}^{\infty} \frac{1}{\ell! (\lambda + \ell)!} \left(\log\left(\frac{z}{2}\right) - \frac{1}{2} \psi(\ell + 1) + \frac{1}{2} \psi(\lambda + \ell + 1)\right).
\]

in (4.12). □

The complex case is analogous. One obtains the following expression.

**Corollary 4.10.** In the complex case the massive Euclidean propagator has an expansion

\[
G_{m,\mathbb{C},D}(x - y) = \frac{1}{(2\pi)^D} \sum_{\ell=0}^{D-2} \frac{(-m^2)^\ell}{\ell! \|x - y\|^{2(D - 2 - \ell)}} \cdot \frac{2^{D-2 - 2\ell} \cdot (D - 2 - \ell)!}{2^\ell \cdot (D - 2 - \ell)!} \\
+ \frac{1}{(2\pi)^D} \sum_{\ell=1}^{\infty} \frac{m^{2(D - 1 + \ell)} \|x - y\|^{2\ell}}{2^{2\ell} \cdot (D - 1 + \ell)!} \cdot \left(\log\left(\frac{m \|x - y\|}{2}\right) - \frac{1}{2} \psi(\ell + 1) + \frac{1}{2} \psi(D + \ell)\right)
\]

This gives a corresponding expansion for the Feynman amplitudes.

**Proposition 4.11.** The Feynman amplitude $\omega_{T,m,R}$ of (3.5) has an expansion

\[
(4.15) \quad \omega_{T,m,R} = \sum_{\ell - (\ell_e) \in \{-\lambda, -\lambda + 1, \ldots, \infty\} \cap E_T} \left(\frac{1}{2\pi}\right)^{(\lambda + 1)|E_T|} \prod_{e \in E_T} B_{T,m,D,\ell_e}^{R} \beta_{T,m,D,\ell_e}^{R} \prod_{v \in V_T} dx_v,
\]

where

\[
(4.16) \quad B_{T,m,D,\ell_e}^{R} := \begin{cases} 
(m^2)^{\lambda - \ell_e \cdot 2^{-\lambda - \ell_e - 1}} \ell_e \in \{-\lambda, \cdots, -1\} \\
(-1)^{\lambda + 1} \frac{m^{2\lambda}}{2^{\lambda + \ell_e \cdot \ell_e + 1}} \ell_e \in \{0, 1, \cdots, \infty\}
\end{cases}
\]

\[
(4.17) \quad \beta_{T,m,D,\ell_e}^{R} := \begin{cases} 
\|x_{a(e)} - x_{t(e)}\|^{2\ell_e} \ell_e \in \{-\lambda, \cdots, -1\} \\
\|x_{a(e)} - x_{t(e)}\|^{2\ell_e + 1} \left(\log\left(\frac{m \|x_{a(e)} - x_{t(e)}\|}{2}\right) - \frac{1}{2} \psi(\ell_e + 1) + \frac{1}{2} \psi(\lambda + \ell_e + 1)\right) \ell_e \in \{0, 1, \cdots, \infty\}
\end{cases}
\]

**Proof.** The statement follows by substituting the Laurent expansion of propagator in Lemma 4.9 in (3.10). □

**Corollary 4.12.** The Feynman amplitude $\omega_{T,m,C}$ of (3.10) has an expansion
\[
\omega_{\Gamma,m,C} = \sum_{\ell=(\ell_e) \in \{-D+1,-D+2,\ldots,\infty\}^{E_{\Gamma}}} \left( \frac{1}{2\pi} \right)^{D[|E_{\Gamma}|]} \prod_{e \in E_{\Gamma}} B_{G_{m,D,\ell_e}}^{C_{\Gamma,m,D,\ell_e}} \beta_{C_{\Gamma,m,D,\ell_e}}^{\ell_e} \int_{x \in V_{\Gamma}} dx_v \wedge d\bar{x}_v,
\]
with \(B_{G_{m,D,\ell_e}}^{C_{\Gamma,m,D,\ell_e}}\) and \(\beta_{C_{\Gamma,m,D,\ell_e}}^{\ell_e}\) respectively as in (4.16) and (4.17), after replacing \(\lambda\) with \(D-1\).

### 4.4.1. Asymptotic expansions of massive propagators

The asymptotic expansion of the Bessel functions also gives a formal expansion for the massive Euclidean propagator in terms of the distance function.

**Lemma 4.13.** The massive Euclidean propagator has an asymptotic expansion of the form

\[
G_{m,R}^{|x-y|} \sim \sqrt{\pi} e^{-m|x-y|} \sum_{\ell=0}^{\infty} \frac{(\lambda,\ell)}{(2\ell+\ell-\lambda+1/2)^{\lambda+1/2}}.
\]

where \(D = 2\lambda + 2\) and

\[
(\lambda,\ell) = \frac{\Gamma(\lambda + \ell + 1/2)}{\ell! \Gamma(\lambda - \ell + 1/2)}.
\]

**Proof.** The modified Bessel function \(K_{\nu}(z)\) has asymptotic expansion, see p. 202 in [25]

\[
K_{\nu}(z) \sim \left( \frac{\pi}{2z} \right)^{1/2} e^{-z} \sum_{\ell=0}^{\infty} \frac{(\nu,\ell)}{(2z)^{\ell}},
\]

with

\[
(\nu,\ell) = \frac{\Gamma(\nu + \ell + 1/2)}{\ell! \Gamma(\nu - \ell + 1/2)}.
\]

\(\square\)

### 4.4.2. Asymptotic expansions of Feynman amplitudes

We obtain a corresponding formal expansion for the Feynman amplitude.

**Proposition 4.14.** The Feynman amplitude \(\omega_{\Gamma,m,R}\) of (3.5) has asymptotic expansion

\[
\omega_{\Gamma,m,R} \sim \sum_{\ell=(\ell_e) \in \mathbb{N}^{E_{\Gamma}}} C_{G_{m,D,\ell_e}}^{R_{\Gamma,m,D,\ell_e}} \omega_{\Gamma,m,D,\ell_e,R},
\]

where

\[
C_{G_{m,D,\ell_e}}^{R_{\Gamma,m,D,\ell_e}} := 2^{-|E_{\Gamma}|/2} \pi^{-(\lambda-\frac{1}{2})|E_{\Gamma}|} \prod_{e \in E_{\Gamma}} \frac{(\lambda,\ell_e)}{2\ell_e+\lambda+1} m_e^{\ell_e-\lambda+1/2}
\]

\[
\omega_{\Gamma,m,D,\ell_e,R} := \prod_{e \in E_{\Gamma}} e^{-m_e \sum_{e} \|x_s(e)-x_t(e)\|} \int_{x \in V_{\Gamma}} dx_v \wedge d\bar{x}_v.
\]

**Proof.** The statement follows by substituting the asymptotic expansion of propagator (4.26) in (3.5). \(\square\)
4.4.3. Complex case. Similarly, we obtain an asymptotic expansion for the massive propagator in the complex case and a corresponding expansion of the Feynman amplitude.

**Corollary 4.15.** The massive Euclidean propagator in the complex case has an asymptotic expansion of the form

\[
G_{m,C,D}(x - y) \sim \sqrt{\frac{\pi}{2}} e^{-m\|x-y\|} \sum_{\ell=0}^{\infty} \frac{(D - 1, \ell)}{2^{\ell}m^{\ell-D+3/2}} \|x - y\|^{-(\ell+D-\frac{1}{2})}.
\]

The Feynman amplitude \(\omega_{\Gamma,m,C}\) of (3.10) has asymptotic expansion

\[
\omega_{\Gamma,m,C} \sim \sum_{\ell=(t_e)\in\mathbb{N}^E_{\Gamma}} C_{T,m,D,\ell}^\mathbb{C} \omega_{\Gamma,m,D,\ell},
\]

where the coefficients \(C_{T,m,D,\ell}^\mathbb{C}\) are as in (4.24), with \(\lambda\) replaced by \(D - 1\), and the forms are given by

\[
\omega_{\Gamma,m,D,\ell} := \prod_{e\in E_{\Gamma}} e^{-m_e \sum_e \|x_{s(e)} - x_{t(e)}\|} \left(\ell_e + D - \frac{1}{2}\right) \bigwedge_{v\in V_{\Gamma}} dx_v \wedge d\bar{x}_v.
\]

**Remark 4.16.** The asymptotic expansion of the Bessel functions exhibits an interesting Stokes phenomenon, see [4]. It would be interesting to interpret this in the setting of Feynman integrals, perhaps along the lines of the role of monodromies in the massive momentum space integrals in [3].

5. Real case: \(x\)-space method and Gegenbauer polynomials

We explain in this section the reason why it is preferable to pass from the real to the complex formulation of Feynman amplitudes. As we already showed in [8], since in the real case the form is not closed, one cannot appeal to cohomological arguments to reformulate the Feynman integral computation, and the only available method is a direct expansion of the Feynman amplitude in orthogonal polynomials and integration on certain chains in the configuration space that lead to nested sums. We have developed the corresponding computation in detail in the massless case in §4 of [8], so we will focus here on how one may be able to set up a similar computational machinery for the massive case. As one already sees in the massless case, the computations quickly become intractable as the graph grows in complexity. By contrast, as we will see in the following section, the complex case leads directly to a much more elegant algebro-geometric formalism for the regularization and renormalization of Feynman amplitudes.

5.1. Gegenbauer polynomials and spherical harmonics. We recall briefly some well known facts about Gegenbauer polynomials and their relation to spherical harmonics, which we need in the following. We refer the reader to [22] for a more detailed treatment.

5.1.1. Generating function an orthogonality. The Gegenbauer polynomials \(C_n^{(\lambda)}(x)\) are determined by the generating function, for \(|t| < 1\),

\[
\frac{1}{(1 - 2tx + t^2)^\lambda} = \sum_{n=0}^{\infty} C_n^{(\lambda)}(x) t^n.
\]

They satisfy the orthogonality relation, for \(\lambda > 1/2\),

\[
\int_{-1}^{1} C_n^{(\lambda)}(x) C_m^{(\lambda)}(x) (1 - x^2)^{\lambda - 1/2} dx = \delta_{n,m} \frac{\pi^{2\lambda - 2} \Gamma(n + 2\lambda)}{n!(n + \lambda)\Gamma(\lambda)^2}.
\]
5.1.2. Zonal and spherical harmonics. Let \( \{ Y_j \} \) be an orthonormal basis of the Hilbert space \( \mathcal{H}_n(S^{D-1}) \) of spherical harmonics on \( S^{D-1} \) of degree \( n \). The zonal spherical harmonics \( Z^{(n)}_{\omega_1}(\omega_2) \) are given by
\[
Z^{(n)}_{\omega_1}(\omega_2) = \dim \mathcal{H}_n(S^{D-1}) \sum_{j=1}^{\dim \mathcal{H}_n(S^{D-1})} Y_j(\omega_1)Y_j(\omega_2).
\]

The Gegenbauer polynomials are related to the zonal spherical harmonics by
\[
C^{(\lambda)}_n(\omega_1 \cdot \omega_2) = c_{D,n} Z^{(n)}_{\omega_1}(\omega_2),
\]
for \( D = 2\lambda + 2 \), with \( \omega_1, \omega_2 \in S^{D-1} \), with the coefficient \( c_{D,n} \) given by
\[
c_{D,n} = \frac{Vol(S^{D-1})(D-2)}{2n + D - 2}.
\]

5.1.3. Expansions in Gegenbauer polynomials. A general procedure for obtaining expansions in Gegenbauer polynomials is described in [2]. For our purposes, it suffices to consider the particular case of the expansion
\[
x^m = 2^{-m} \Gamma(\lambda)m! \sum_{k=0}^{[m/2]} \frac{(\lambda + m - 2k)}{k!(\lambda + m + 1 - k)} C^{(\lambda)}_{m-2k}(x).
\]
for a non-negative integer \( m \) and for \( \lambda \geq 1/2 \).

5.1.4. Chebyshev polynomials of the first kind. The Chebyshev polynomials of the first kind have generating function
\[
-\frac{1}{2} \log(1 - 2tx + t^2) = \sum_{n=1}^{\infty} T_n(x) \frac{t^n}{n},
\]
They are expressible as a limit of Gegenbauer polynomials,
\[
T_n(x) = \frac{n}{2} \lim_{\lambda \to 0} \frac{C^{(\lambda)}_n(x)}{\lambda}.
\]
They can also be written as combinations of Gegenbauer polynomials of a fixed weight \( \lambda \): the expression as a function of Gegenbauer polynomials can be obtained inductively, using the recurrence relation \( T_0(x) = 1, T_1(x) = x \), and \( T_{n+1}(x) = 2xT_n(x) - T_{n-1}(x) \) for the Chebyshev polynomials, together with [5.6]. We denote by \( c^{\lambda}_{n,m} \) the resulting coefficients,
\[
T_n(x) = \sum_{m=0}^{n} c^{\lambda}_{n,m} C^{(\lambda)}_m(x).
\]

5.2. Expansion of Feynman amplitudes in Gegenbauer polynomials. We described in §4 of [8] how to expand the real Euclidean massless Feynman amplitudes in Gegenbauer polynomials. Here we extend the procedure to the massive case. We use the expansion of the massive Feynman amplitudes given in Proposition 4.11.
Theorem 5.1. The forms \((2\pi)^{-(\lambda+1)}B_{\Gamma,m,D,\ell,e}^R, \beta_{\Gamma,m,D,\ell,e}^R\) in the expansion (4.15) of the real massive Euclidean Feynman amplitude have a series expansion in Gegenbauer polynomials, either of the form

\[
\frac{B_{\Gamma,m,D,\ell,e}^R}{(2\pi)^{(\lambda+1)}} \beta_{\Gamma,m,D,\ell,e}^R = \rho_e^{-2\ell e} \sum_{n_e=0}^{\infty} \left( \frac{r_e}{\rho_e} \right)^{n_e} \sum_{k_e=0}^{[n_e/2]} \sum_{j_e=0}^{[(n_e-2k_e)/2]} d_{n_e,k_e,j_e}^\lambda C_{n_e-2(k_e+j_e)}^{(\lambda)} \left( \omega_{s(e)} \cdot \omega_{t(e)} \right)
\]

or of the form

\[
\frac{B_{\Gamma,m,D,\ell,e}^R}{(2\pi)^{(\lambda+1)}} \beta_{\Gamma,m,D,\ell,e}^R = \rho_e^{2\ell e} \log \left( \frac{m_e}{2\rho_e} \right) \sum_{n_e=0}^{\infty} \sum_{k_e=0}^{[n_e/2]} \sum_{p_e=0}^{\infty} \frac{1}{2^{p_e}} \left( \frac{r_e}{\rho_e} \right)^{n_e+p_e} \sum_{s_e=0}^{p_e} \sum_{r_e=0}^{[(p_e+s_e)/2]} \alpha_{r_e,p_e,m_e}^\lambda \beta_{r_e,j_e,p_e,s_e}^{(\ell_e)} C_{p_e+s_e-2(r_e+j_e)}^{(\lambda)} \left( \omega_{s(e)} \cdot \omega_{t(e)} \right).
\]

In the case where \(D = 2\lambda + 2\) with \(\lambda\) an integer, the expansion has coefficients \(a^\lambda, b^\lambda, c^\lambda, d^\lambda, f^\lambda, \alpha^\lambda, \) and \(\beta^\lambda\) in

\[Q[m_e, \log(m_e), -1, \gamma, \log(2)],\]

where \(\gamma\) is the Euler-Mascheroni constant. In the case with \(D = 2\lambda + 2\) with \(\lambda\) a half-integer, the expansion has coefficients in

\[Q[m_e, \log(m_e), -\frac{1}{2}, \gamma, \log(2)].\]

Proof. Consider the forms \(\beta_{\Gamma,m,D,\ell,e}^R\), given as in (4.17). In the case

\[\beta_{\Gamma,m,D,\ell,e}^R = \|x_{s(e)} - x_{t(e)}\|^{-2\ell e}\]

we have an expansion

\[\|x_{s(e)} - x_{t(e)}\|^{-2\ell e} = \rho_e^{-2\ell e} \left( 1 + \left( \frac{r_e}{\rho_e} \right)^2 - \frac{r_e}{\rho_e} \omega_{s(e)} \cdot \omega_{t(e)} \right)^{-\ell e} = \rho_e^{-2\ell e} \sum_{n_e=0}^{\infty} \frac{C_{n_e}^{(\ell_e)}(\omega_{s(e)} \cdot \omega_{t(e)})}{\left( \frac{r_e}{\rho_e} \right)^{n_e}}.\]

This is already an expansion in Gegenbauer polynomials, but with varying weight \(\ell_e\). We need instead an expansion in the Gegenbauer polynomials \(C_n^{(\lambda)}\) with the fixed \(\lambda\) related to the spacetime dimension by \(D = 2\lambda + 2\), which are the polynomials related to the spherical harmonics on \(S^{D-1}\).

We can turn the expansion above into an expansion in Gegenbauer polynomials \(C_n^{(\lambda)}\), using the explicit expression for the Gegenbauer polynomials

\[(5.10) \quad C_n^{(\lambda)}(x) = \sum_{k=0}^{[n/2]} (-1)^k \frac{\Gamma(\lambda + n - k)}{k! (n - 2k)! \Gamma(\lambda)} (2x)^{n-2k},\]

together with (5.6), which gives

\[(5.11) \quad C_{n_e}^{(\ell_e)}(x) = \sum_{k=0}^{[n_e/2]} (-1)^k \frac{\Gamma(\ell_e + n_e - k) \Gamma(\lambda)}{k! \Gamma(\ell_e)} \frac{[(n_e-2k)/2]}{\Gamma(\lambda + n_e - 2(k + j))} C_{n_e-2(k+j)}^{(\lambda)}(x).
\]

Thus, we obtain

\[\|x_{s(e)} - x_{t(e)}\|^{-2\ell} = \]
\[
\rho_e^{-2\ell_e} \sum_{n_e=0}^\infty \left( \frac{r_e}{\rho_e} \right)^{[n_e/2]} \sum_{k=0}^{[n_e/2]} (-1)^k \Gamma(\ell_e + n_e - k) \frac{\Gamma(\lambda)}{k! \Gamma(\ell_e)} \sum_{j=0}^{[n_e-2k]/2} \frac{(\lambda + n_e - 2(k + j))}{j! \Gamma(\lambda + n_e - 2k + 1 - j)} C_n^{(\lambda)} \omega_{n-e-2(k+j)}(\omega_{s(e)} \cdot \omega_{t(e)}).
\]

In the case where
\[
\beta^{\mathbb{R}}_{T,m,D,e} = \|x_{s(e)} - x_{t(e)}\|^{2\ell} \log(\frac{m_e}{2}) \log(\frac{\|x_{s(e)} - x_{t(e)}\|}{2}) - \frac{\psi(\ell_e + 1 + \psi(\lambda + \ell_e + 1)}{2}
\]

we take the Taylor expansion
\[
\rho_e^{2\ell_e} \left( 1 - 2 \frac{r_e}{\rho_e} x + \left( \frac{r_e}{\rho_e} \right)^2 \right)^{\ell_e} =
\]
\[
\rho_e^{2\ell_e} \sum_{n=0}^{\ell_e} \frac{\Gamma(\ell_e + 1)}{\Gamma(\ell_e - n)} (-1)^n \left( \frac{r_e}{\rho_e} \right)^n (1 + \left( \frac{r_e}{\rho_e} \right)^2)^{\ell_e - n} \sum_{k=0}^{n/2} \frac{(\lambda + n - 2k)}{k! \Gamma(\lambda + n + 1 - k)} C_{n-2k}^{(\lambda)} \omega_{n-e-2(k+j)}(\omega_{s(e)} \cdot \omega_{t(e)}).
\]

We also use the expansion into Chebyshev polynomials and the expression of Chebyshev polynomials in terms of Gegenbauer polynomials. This gives
\[
\log(\frac{m_e}{2}) \log(\frac{\|x_{s(e)} - x_{t(e)}\|}{2}) = \log(m_e) - \log(2) + \log(\|x_{s(e)} - x_{t(e)}\|)
\]

where the first terms contribute coefficients in \(\mathbb{Q}[\log(m_e), \log(2)]\) to the expansion and the last term is expanded as
\[
\log(\|x_{s(e)} - x_{t(e)}\|) = \log(\rho_e) + \log((1 + \left( \frac{r_e}{\rho_e} \right)^2 - 2 \frac{r_e}{\rho_e} \omega_{s(e)} \cdot \omega_{t(e)})^{1/2}) =
\]
\[
\log(\rho_e) - \sum_{n=1}^{\infty} \frac{1}{n} \left( \frac{r_e}{\rho_e} \right)^n \sum_{m=0}^{n} C_{n,m}^{(\lambda)} \omega_{s(e)} \cdot \omega_{t(e)}).
\]

As one can see from [5.6], the coefficients \(C_{n,m}^{(\lambda)}\) are in \(\mathbb{Q}\) when \(\lambda\) is an integer. When \(\lambda\) is a half-integer, the coefficients contain values of the Gamma function at half integers. By the duplication formula \(\Gamma(z)\Gamma(z + 1/2) = 2^{1-2z} \sqrt{\pi} \Gamma(2z)\), the values of \(\Gamma(z)\) at half integers are in \(\mathbb{Q}[\sqrt{\pi}]\), hence the \(C_{n,m}^{(\lambda)}\) are in \(\mathbb{Q}[\sqrt{\pi}]\).

Finally, observe that the products of Gegenbauer polynomials that we obtain when we combine the expansion of \(\|x_{s(e)} - x_{t(e)}\|^{2\ell}\) with that of \(\log(\|x_{s(e)} - x_{t(e)}\|)\) in the product term, can in turn be expressed in terms of a sum of Gegenbauer polynomials. In fact, using [5.10] we have
\[
C_n^{(\lambda)}(x) C_m^{(\lambda)}(x) = \sum_{k=0}^{[n/2]} \sum_{j=0}^{[m/2]} (-1)^{j+k} \Gamma(\lambda + n - k) \Gamma(\lambda + m - j) \frac{k! j! (n-2k)! (m-2j)! \Gamma(\lambda)^2 (2x)^{n+m-2(k+j)}}{k! j! (n-2k)! (m-2j)!}
\]

Setting
\[
\alpha^{\lambda}_{r,n,m} = \frac{(-1)^{2\ell + n - 2r}}{\Gamma(\lambda)^2} 2^{-(n+m-2r)} \Gamma(\lambda)(n + m - 2r) \sum_{j+k=r} \Gamma(\lambda + n - k) \Gamma(\lambda + m - j) \frac{k! j! (n-2k)! (m-2j)!}{k! j! (n-2k)! (m-2j)!}
\]
we obtain
\[
C_n^{(\lambda)}(x) C_m^{(\lambda)}(x) = \sum_{r=0}^{[(m+n)/2]} \sum_{k=0}^{[(n+m-2r)/2]} \beta_{r,k,n,m}^{\lambda} C_{n+m-2(r+k)}^{(\lambda)}(x)
\]

When \(\lambda\) is an integer, the values \(\Gamma(\lambda+n), n \in \mathbb{N}\) of the Gamma function are integers, while the values of the Digamma function \(\psi(\lambda+n), n \in \mathbb{N}\) are of the form \(H_{\lambda+n-1} - \gamma\) when \(\lambda\) is an integer, where \(H_n\) is the \(n\)-th harmonic number and \(\gamma\) is the Euler–Mascheroni constant. Thus, the coefficients of the expansion of the Feynman amplitude obtained in this way are in \(\mathbb{Q}[m_e, \log(m_e), (2\pi)^{-1}, \gamma]\). When \(\lambda\) is a half-integer, we also have coefficients that have, in the numerator or denominator, values at half-integers of the Gamma function, and which are therefore in \(\mathbb{Q}[\sqrt{\pi}^{\pm 1}]\). Moreover, when \(\lambda\) is a half-integer, there are also coefficients in the expansion that contain values of the Digamma function at half-integers and these are of the form \(\psi(n + \frac{1}{2}) = -\gamma - 2\log(2) + \sum_{k=1}^{n} \frac{\pi}{2\Gamma(k)}\), in \(\mathbb{Q}[\gamma, \log(2)]\). Therefore, the coefficients must be in \(\mathbb{Q}[m_e, \log(m_e), \sqrt{\pi}^{\pm 1}, \gamma, \log(2)]\).

We can then apply the same computational method described in §4 of [8] to the terms of this expansion, using the decomposition of the chain of integration \(\sigma_{\Gamma} : \mathbb{R} \to \mathbb{R}\) into the chain described in [8] associated to acyclic orientations of the graph. That leads to subdividing the integrals into an angular part, that can be expressed in terms of integrals of products of Gegenbauer polynomials at the star of half-edges around each vertex of the graph, followed by an integrals in the radial coordinates of the resulting function depending on the matching pairs of half edges. The method follows the same pattern described in §4 of [8] and we do not discuss it detail here. Clearly, although in principle explicitly computable by this method, computations of Feynman integrals performed with this \(x\)-space method become very quickly intractable as the complexity of the graph increases. Thus, we turn instead to a different method, based on using the complexified form of the Feynman amplitude (3.10) and algebro-geometric techniques using the wonderful compactification spaces \(X_H[n]\).

6. Complex case: forms with logarithmic poles and Rota–Baxter renormalization

In this section we consider the complex case of Feynman amplitudes, as introduced in §3.2.3 in the general framework of algebraic renormalization. In this section we consider the Rota–Baxter algebra of Feynman graphs and a Rota–Baxter algebra, which we construct using complexes of differential forms with logarithmic poles on the wonderful compactifications \(X_H[n]\). The renormalization procedure for Feynman amplitudes in this setting consists of the following steps:

1. **Pullback** of the form \(\omega_{\Gamma,C}\) to \(X_H[n]\), under the map \(\pi_n : X_H[n] \to X^n\) of the iterated blowup construction.

2. **Regularization** of \(\pi_n^*(\omega_{\Gamma,C})\) by replacing this closed form with a cohomologous form \(\eta_{\Gamma}\) on \(X_H[n]\) with logarithmic poles along the boundary divisor \(D_n\) of the compactification \(X_H[n]\).

3. **Construction** of a Rota–Baxter algebra \(\mathcal{R}\) of forms on \(X_H[n]\) with logarithmic poles along \(D_n\) with a linear operator of “polar subtraction” defined in terms of Poincaré residues.

4. **An algebraic Feynman rule** given by a commutative algebra homomorphism \(\phi : H \to \mathcal{R}\) from the Connes-Kreimer Hopf algebra of Feynman graphs \(H\) to the Rota–Baxter algebra, given by \(\phi(\Gamma) = \eta_{\Gamma}\).

5. **Renormalization** of the Feynman amplitude by Birkhoff factorization of this algebraic Feynman rule, using the coproduct and antipode of \(H\) and the Rota–Baxter operator of \(\mathcal{R}\).

6. **Evaluation** of the renormalized Feynman amplitude on the chain \(\tilde{\sigma}_{\Gamma,C} = \pi_n^{-1}(\sigma_{\Gamma,C})\).
We also show that, in this setting we obtain the beta function of the renormalization group as an element in the Lie algebra $\mathfrak{g}(\mathcal{R})$ of the group $G(\mathcal{R}) = \text{Hom}_{\text{Alg}}(\mathfrak{h}, \mathcal{R})$.

6.1. **Algebraic renormalization.** A general algebraic framework for renormalization was formulated in terms of Rota–Baxter algebras and Hopf algebras in [13], following the analysis of the BPHZ renormalization procedure carried out in [9]. We recall here the setting, as given in [13].

### 6.1.1. Rota–Baxter algebras.

A Rota–Baxter algebra of weight $\lambda = -1$ is a commutative unital algebra $R$ with a linear operator $T : R \to R$ satisfying the Rota–Baxter relation

$$T(x)T(y) = T(xT(y)) + T(T(x)y) + \lambda T(xy).$$

The main example that corresponds to the use of dimensional regularization of Feynman integrals in [9] is the algebra of Laurent series with $T$ given by the projection onto the polar part.

The Rota–Baxter operator $T$ determines a splitting of $R$ into a piece $R_+ = (1 - T)R$ and a piece $R_-$ which is the unitization of $TR$. Both the $R_\pm$ are algebras, due to the Rota–Baxter relation (6.1).

### 6.1.2. Connes–Kreimer Hopf algebra of Feynman graphs.

We consider here the usual Connes–Kreimer Hopf algebra $H$ of Feynman graphs [9]. This is a free commutative algebra with generators the 1PI Feynman graphs $\Gamma$ of the theory, graded by the number of internal lines, with

$$\deg(\Gamma_1 \cdots \Gamma_n) = \sum_i \deg(\Gamma_i), \quad \deg(1) = 0,$$

and with coproduct

$$\Delta(\Gamma) = \Gamma \otimes 1 + 1 \otimes \Gamma + \sum_{\gamma \in V(\Gamma)} \gamma \otimes \Gamma / \gamma$$

where the sum is over all subgraphs $\gamma$ such that the quotient $\Gamma / \gamma$ is still a 1PI Feynman graph of the same theory. The antipode is defined inductively by

$$S(X) = -X - \sum S(X')X''$$

for $\Delta(X) = X \otimes 1 + 1 \otimes X + \sum X' \otimes X''$, with the $X'$ and $X''$ of lower degrees.

### 6.1.3. Algebraic Feynman rules.

In this algebraic setting a Feynman rule is simply a morphism of commutative algebras from a commutative Hopf algebra $\mathfrak{h}$ to a Rota–Baxter algebra $\mathcal{R}$ of weight $-1$,

$$\phi \in \text{Hom}_{\text{Alg}}(\mathfrak{h}, \mathcal{R}).$$

Notice that the morphism $\phi$ only preserves the commutative algebra structure and has no required compatibility with the Hopf and the Rota–Baxter structures.

### 6.1.4. Birkhoff factorization.

As was originally shown by Connes and Kreimer in [9] and generalized to the Rota–Baxter setting in [13], given an algebraic Feynman rule $\phi \in \text{Hom}_{\text{Alg}}(\mathfrak{h}, \mathcal{R})$ as above, the Hopf structure on $\mathfrak{h}$ together with the Rota–Baxter structure on $\mathcal{R}$ determine a Birkhoff factorization

$$\phi = (\phi_- \circ S) \ast \phi_+$$

into commutative algebra homomorphisms

$$\phi_\pm \in \text{Hom}_{\text{Alg}}(\mathfrak{h}, \mathcal{R}_\pm).$$
Lemma 6.1. Let \( D \) be a normal crossings divisor in a smooth variety \( Y \) and let \( \Omega^k_Y(\log D) \) denote the sheaf of differential forms on \( Y \) with logarithmic poles along \( D \), as in [11]. Let \( \{ D_{j} \}_{j \in I} \) denote the components of \( D \) and, for \( J \subseteq I \), \( J = \{ j_1, \ldots, j_r \} \), let \( D_J = D_{j_1} \cap \cdots \cap D_{j_r} \) denote the intersection.

Given a form \( \eta \in \Omega^k_Y(\log D) \) and a nonempty intersection \( D_J \), there is an iterated Poincaré residue \( \text{Res}_{D_J}(\eta) \), which is a form in \( \Omega^{k-r}_{D_J} \). The pairing of \( \text{Res}_{D_J}(\eta) \) with a \((k-r)\)-cycle \( \Sigma \) in \( D_J \) is given by

\[
\int_{\Sigma} \text{Res}_{D_J}(\eta) = \frac{1}{(2\pi i)^r} \int_{\mathcal{L}_{D_J}(\Sigma)} \eta,
\]

where \( \mathcal{L}_{D_J}(\Sigma) \) is the iterated Leray coboundary of \( \Sigma \), which is a \( k \)-cycle in \( Y \) given by a \( T^r \)-torus bundle over \( \Sigma \) obtained as composition of the Leray coboundary maps

\[
\mathcal{L}_{D_J}(\Sigma) = \mathcal{L}_{D_{j_1}} \circ \cdots \circ \mathcal{L}_{D_{j_k}}(\Sigma)
\]

of the components of \( D_J \), see Proposition 5.14 of [8].

**Lemma 6.1.** Let \( f_1, \ldots, f_n \) be local parameters such that the components \( D_j \) of the divisor \( D_n \) are defined by \( f_j = 0 \). Given \( \eta \in \Omega_Y^m(\log D) \), set

\[
T(\eta) = \sum_{j=1}^{n} \frac{df_j}{f_j} \wedge \text{Res}_{D_j}(\eta).
\]

Then, for any \( D_J = D_{j_1} \cup \cdots \cup D_{j_r} \) for a subset \( J \subseteq I \) of components of \( D_n \), the form \( T(\eta) \) satisfies

\[
\text{Res}_{D_J}(T(\eta)) = \text{Res}_{D_J}(\eta).
\]

**Proof.** First observe that for the ordinary Poincaré residue for a single component \( D_j \) we have \( \text{Res}_j(T(\eta)) = \text{Res}_j(\eta) \) by construction. Then, using (6.4) and (6.5) we see that the iterated residues also agree,

\[
\langle \text{Res}_{D_J}(T(\eta)), \Sigma \rangle = \langle T(\eta), \mathcal{L}_{D_J}(\Sigma) \rangle
= \langle \text{Res}_{j_1}(T(\eta)), \mathcal{L}_{j_2} \cdots \mathcal{L}_{j_k}(\Sigma) \rangle
= \langle \text{Res}_{j_1}(\eta), \mathcal{L}_{j_2} \cdots \mathcal{L}_{j_k}(\Sigma) \rangle = \langle \text{Res}_{D_J}(\eta), \Sigma \rangle.
\]

\[\square\]
Given a form \( \eta \in \Omega_y^\bullet(\log D) \), we define the polar subtraction \( \eta_D \) to be
\[
\eta_D := \eta - T(\eta).
\]
We will show below that this operation of polar subtraction has the same formal properties of other polar subtraction methods used in quantum field theory (such as DimReg with MS), and can be treated similarly to obtain a good notion of renormalized Feynman integrals in configuration spaces. The result of the polar subtraction is a form \( \eta_D \) for which all the iterated Poincaré residues vanish.

**Corollary 6.2.** The form \( \eta_D \) has vanishing Poincaré residues, \( \text{Res}_{D_J}(\eta_D) = 0 \) for all \( J \subset I \).

**Proof.** This is an immediate consequence of Lemma 6.1. \(\square\)

### 6.2.2. Rota–Baxter algebra of the wonderful compactifications

We now consider the wonderful compactifications \( X_{H[n]} \) described in Section 2 above, with the boundary divisor
\[
D_n = X_{H[n]} \setminus \left( (X \setminus H)^n \setminus \bigcup_{\{i,j\} \subset n} \Delta_{\{i,j\}} \right).
\]

Let \( \Omega^\text{even}_Y(\log D) \) be the commutative algebra of even differential forms on a variety \( Y \) with logarithmic poles along a normal crossings divisor \( D \).

Define the linear operator \( T : \Omega^\text{even}_Y(\log D) \to \Omega^\text{even}_Y(\log D) \) given as in (6.6) by
\[
T : \eta \mapsto T(\eta) = \sum_{j=1}^n \frac{df_j}{f_j} \wedge \text{Res}_{D_j}(\eta).
\]

**Theorem 6.3.** The pair \( (R, T) \), with \( R = \Omega^\text{even}_Y(\log D) \) and with \( T \) as in (6.10), is a Rota–Baxter algebra of weight \(-1\).

**Proof.** First observe that the operator \( T \) is idempotent: in fact, by Lemma 6.1, we have
\[
T(T(\eta)) = \sum_{j=1}^n \frac{df_j}{f_j} \wedge \text{Res}_{D_j}(T(\eta)) = T(\eta).
\]

We need to check that the Rota–Baxter relation (6.1) is satisfied. We have
\[
T(\eta) \wedge T(\xi) = \sum_{j<k} \frac{df_j}{f_j} \wedge \text{Res}_{D_j}(\eta) \wedge \sum_{k} \frac{df_k}{f_k} \wedge \text{Res}_{D_k}(\xi),
\]
which can equivalently be written as
\[
T(\eta) \wedge T(\xi) = \sum_{j<k} \frac{df_j}{f_j} \wedge \frac{df_k}{f_k} \wedge (\text{Res}_{D_k}(\eta) \wedge \text{Res}_{D_j}(\xi) - \text{Res}_{D_j}(\eta) \wedge \text{Res}_{D_k}(\xi)).
\]

For the term \( T(\eta \wedge \xi) \) we obtain
\[
T(\eta \wedge \xi) = \sum_{j} \frac{df_j}{f_j} \wedge \text{Res}_{D_j}(\eta \wedge \xi)
= \sum_{j} \frac{df_j}{f_j} \wedge \text{Res}_{D_j}((T(\eta) + (\eta - T(\eta))) \wedge (T(\xi) + (\xi - T(\xi))))
= \sum_{j} \frac{df_j}{f_j} \wedge \text{Res}_{D_j}(T(\eta) \wedge T(\xi))
\]

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We first check that (6.15) determines a Rota–Baxter structure of weight $\eta$ and then inductively that it determines a linear operator $T$ on the Rota–Baxter algebra where all values of $\eta$ combine the Rota–Baxter algebras $\Omega(T)$ as above. Let $R \in \Omega(T)$.

Theorem 6.4. Consider the commutative algebra

$$\mathcal{R} = \bigwedge_n \Omega_{\mathcal{X}_\mathcal{H}[n]}^{\text{even}}(\log D_n),$$

as above. Let $T : \mathcal{R} \to \mathcal{R}$ be the linear operator defined by setting

$$T(\eta_1 \wedge \eta_2) = T_{n_1}(\eta_1) \wedge \eta_2 + \eta_1 \wedge T_{n_2}(\eta_2) - T_{n_1}(\eta_1) \wedge T_{n_2}(\eta_2),$$

for $\eta_i \in \Omega_{\mathcal{X}_\mathcal{H}[n_i]}^{\text{even}}(\log D_{n_i})$. Then $(\mathcal{R}, T)$ is a Rota–Baxter algebra of weight $-1$.

Proof. We first check that (6.15) determines a Rota–Baxter structure of weight $-1$ on

$$\Omega_{\mathcal{X}_\mathcal{H}[n_1]}^{\text{even}}(\log D_{n_1}) \wedge \Omega_{\mathcal{X}_\mathcal{H}[n_2]}^{\text{even}}(\log D_{n_2})$$

and then inductively that it determines a linear operator $T$ of $\mathcal{R}$ that also satisfies the Rota–Baxter relation of weight $-1$. By (6.15) we have

$$T(T_{n_1}(\eta_1) \wedge \eta_2) = T_{n_1}(\eta_1) \wedge \eta_2 + T_{n_1}(\eta_1) \wedge T_{n_2}(\eta_2) - T_{n_1}(\eta_1) \wedge T_{n_2}(\eta_2) = T_{n_1}(\eta_1) \wedge T_{n_2}(\eta_2),$$

$$T(\eta_1 \wedge T_{n_2}(\eta_2)) = T_{n_1}(\eta_1) \wedge T_{n_2}(\eta_2) + \eta_1 \wedge T_{n_2}(\eta_2) - T_{n_1}(\eta_1) \wedge T_{n_2}(\eta_2) = \eta_1 \wedge T_{n_2}(\eta_2),$$

and

$$T^2(\eta) = T(T(\eta)).$$

The last term has a wedge of two forms with no residues on all the $D_j$; in the two intermediate terms the forms $\eta_D = \eta - T(\eta)$ and $\xi_D = \xi - T(\xi)$ have no residue, hence using (6.11) we obtain, respectively, $T(\eta) \wedge (\xi - T(\xi))$ and $(\eta - T(\eta)) \wedge T(\xi)$. Since the forms $\eta$ and $\xi$ are even, commuting them with a 1-form $df/\eta$ does not give rise to a non-trivial sign. Finally, the first term can be computed using the explicit form of $T(\eta)$ and $T(\xi)$ and gives again the expression (6.12). Thus, we have

$$(6.13) \quad T(\eta \wedge \xi) = T(\eta)T(\xi) + T(\eta) \wedge (\xi - T(\xi)) + (\eta - T(\eta)) \wedge T(\xi).$$

We then compute the remaining terms $T(\eta \wedge T(\xi))$ and $T(T(\eta) \wedge \xi)$. Applying (6.13) to these cases and using (6.11) we obtain

$$T(\eta \wedge T(\xi)) = T(\eta) \wedge (\xi - T(\xi)) + (\eta - T(\eta)) \wedge T(\xi),$$

$$T(T(\eta) \wedge \xi) = T(\eta) \wedge (\xi - T(\xi)) + T(\eta) \wedge (\xi - T(\xi)).$$

so that we obtain the Rota–Baxter relation of weight $-1$ as stated,

$$T(\eta \wedge T(\xi)) + T(T(\eta) \wedge \xi) - T(\eta) \wedge T(\xi) = T(\eta \wedge \xi).$$

□
hence the Rota–Baxter relation is satisfied. One easily checks also that the operator $T$ defined in this way still satisfies $T^2 = T$. Now assume inductively that, on all $k$-fold wedge products
\[
\Omega^\text{even}_{g([n_1])} (\log D_{n_1}) \wedge \cdots \wedge \Omega^\text{even}_{g([n_k])} (\log D_{n_k})
\]
we have constructed $T$ satisfying the Rota–Baxter relation of weight $-1$, with $T^2 = T$. Then we can extend it to a $(k + 1)$-fold wedge product by setting
\[
T(\eta_1 \wedge \cdots \wedge \eta_k \wedge \eta_{k+1}) = T(\eta_1 \wedge \cdots \wedge \eta_k) \wedge \eta_{k+1}
\]
+ $\eta_1 \wedge \cdots \wedge \eta_k \wedge T_{n_{k+1}}(\eta_{k+1}) - T(\eta_1 \wedge \cdots \wedge \eta_k) \wedge T_{n_{k+1}}(\eta_{k+1}).$

The same argument above shows that it still satisfied the Rota–Baxter relation and one sees in the same way that $T^2 = T$. □

We refer to the Rota–Baxter algebra obtained as in Theorem 6.4 as the Rota–Baxter algebra of configuration spaces.

6.3. Algebraic renormalization of Feynman integrals in configuration space.

6.3.1. The algebraic Feynman rule. Let $\omega_{\Gamma, C}$ be the massless Feynman amplitude for a Feynman graph $\Gamma$ in the complex case, defined as in (3.8), as a smooth closed form in $\Omega^{2D_n}(F(\mathcal{X} \setminus \mathcal{H}, n))$. As in [3], we then consider the pullback $\pi_n^* (\omega_{\Gamma, C})$ under the iterated blowup map $\pi_n : \mathcal{X}_H[n] \to \mathcal{X}^n$. The pullback form is a smooth closed form in $\Omega^{2D_n}(\mathcal{X}_H[n] \setminus D_n)$, where $D_n$ is the boundary divisor (6.9).

It is well known that the de Rham cohomology of a smooth quasi-projective varieties computed using algebraic differential forms, [15]. Moreover, if the variety is the complement of a normal crossings divisor one can use algebraic differential forms with logarithmic poles along the divisors, [11], so that we have
\[
H^*(\mathcal{X}_H[n] \setminus D_n) \cong H^* (\mathcal{X}_H[n], \Omega^* \mathcal{X}_H[n](\log (D_n))).
\]

Thus, the smooth closed differential form $\pi_n^* (\omega_{\Gamma, C})$ is cohomologous to an algebraic differential form $\eta_{\Gamma}$ with logarithmic poles along $D_n$,
\[
[\pi_n^* (\omega_{\Gamma, C})] = [\eta_{\Gamma}] \in H^{2D_n} (\mathcal{X}_H[n] \setminus D_n).
\]

Let $\mathcal{R}$ be the Rota–Baxter algebra of configuration spaces as in Definition 6.4. We define the algebraic Feynman rule $\phi \in \text{Hom}_{\mathcal{R}} (\mathcal{S}, \mathcal{R})$, for the algebraic renormalization of Feynman integrals in configuration spaces, by setting
\[
(6.16) \quad \phi (\Gamma_1 \cdots \Gamma_k) = \eta_{\Gamma_1} \wedge \cdots \wedge \eta_{\Gamma_k},
\]
where for Feynman graph $\Gamma$ the form $\eta_{\Gamma} \in \Omega^{2D_n}(\mathcal{X}_H[n], D_n)$, with $n = |E_{\Gamma}|$, is the algebraic differential form with logarithmic poles along the divisor $D_n$ of (6.9), obtained as above.

6.3.2. The Birkhoff factorization. The resulting Birkhoff factorization of the Feynman rule (6.16) determines $\phi_-$ and $\phi_+$ inductively as
\[
\phi_- (\Gamma) = -T (\eta_{\Gamma} + \sum_{\gamma \in \Gamma} \phi_- (\gamma) \wedge \eta_{\Gamma/\gamma})
\]
\[
\phi_+ (\Gamma) = (1 - T) (\eta_{\Gamma} + \sum_{\gamma \in \Gamma} \phi_- (\gamma) \wedge \eta_{\Gamma/\gamma}) = \eta_{\Gamma, D} + \sum_{\gamma \in \Gamma} (\phi_- (\gamma) \wedge \eta_{\Gamma/\gamma}) D,
\]
with $\eta_D$ the polar subtraction defined as in (6.8), and similarly for the other terms. By analogy to the usual setting of BPHZ renormalization (see [9]), we refer to the expression $\eta_{\Gamma} + \sum_{\gamma \in \Gamma} \phi_- (\gamma) \wedge \eta_{\Gamma/\gamma}$ as the preparation of the form $\eta_{\Gamma} = \phi (\Gamma)$, to the $\phi_- (\Gamma)$ as the counterterms, and to the $\phi_+ (\Gamma)$ as the renormalized Feynman amplitude in configuration space.
6.3.3. Renormalized Feynman integrals. One then obtains, as renormalized Feynman integral in configuration space the integral

\[(6.17) \int_{\tilde{\Gamma},\mathcal{C}} \eta_{\tilde{\Gamma},\mathcal{D}} + \sum_{\gamma \in \Gamma} (\phi_{-}\gamma) \wedge \eta_{\gamma/\gamma} \mathcal{D},\]

where \( \tilde{\sigma}_{\Gamma,\mathcal{C}} = \pi_{n}^{-1}(\sigma_{\Gamma,\mathcal{C}}) = \mathbb{P}_{H}^{\mathcal{D}}[n] \times \{ y \} \) is the preimage in \( \mathcal{X}_{H}[n] \) of the chain of integration \( \sigma_{\Gamma,\mathcal{C}} \) of the unrenormalized Feynman amplitude \((3.13)\) in the complex case. The form \( \eta_{\tilde{\Gamma},\mathcal{D}} + \sum_{\gamma \in \Gamma} (\phi_{-}\gamma) \wedge \eta_{\gamma/\gamma} \) is now has vanishing residues on all the divisors \( \mathcal{D}_{J} \). Hence, the integral \((6.17)\) is free of divergences.

**Lemma 6.5.** If the motives \( m(\mathcal{X}) \) and \( m(\mathcal{H}) \) of the varieties \( \mathcal{X} \) and \( \mathcal{H} \) are mixed Tate, then the motive \( m(\mathcal{X}_{H}[n]) \) of the configuration space is also mixed Tate, and so are the motives \( m(\mathcal{D}) \) of all the unions and intersections \( \mathcal{D} \) of components of the boundary divisor \( \mathcal{D}_{n} \). In particular the mixed motive \( m(\mathcal{X}_{H}[n], \mathcal{D}_{n}) \) is also mixed Tate.

**Proof.** Due to Voevodsky (Proposition 3.5.3 in \([24]\), the motive \( m(\tilde{\mathcal{Y}}) \) of the blow-up of a smooth scheme \( \mathcal{Y} \) along a smooth closed subscheme \( \mathcal{V} \subset \mathcal{Y} \) is canonically isomorphic to

\[m(\mathcal{Y}) \oplus \bigoplus_{k=1}^{\text{codim}_{\mathcal{Y}}(\mathcal{V})-1} m(\mathcal{V})(k)[2k]\]

where \([-]\) denotes the degree shift in the triangulated category of mixed motives, while \((-)\) is the Tate twist. The result for \( m(\mathcal{X}_{H}[n]) \) simply follows from the repeated application of this blowup formula to the iterated blowup construction of \( \mathcal{X}_{H}[n] \). The case of the motives \( m(\mathcal{D}) \) follows from the explicit stratification of \( \mathcal{X}_{H}[n] \), see \([16]\), where the intersections of the components are described explicitly in terms of other configuration spaces whose motive can be computed by the same method.

Knowing that all the intersections are still mixed Tate motives then implies inductively that the unions also are, by repeatedly using the distinguished triangles in the triangulated category of mixed motives.

Thus, we can formulate the renormalized Feynman integral as a period, in the sense of \([17]\), in the following way.

**Corollary 6.6.** The renormalized Feynman integral \((6.17)\) is a period of a mixed Tate motive over the field of definition of the algebraic form \( \phi_{+}(\eta_{\Gamma}) \).

**Proof.** The motive of the wonderful compactification \( \mathcal{X}_{H}[n] \) is mixed Tate whenever the motive of \( \mathcal{X} \) is. This follows the same lines as the argument used in \([7]\) for the graph configuration spaces, using the description of \( \mathcal{X}_{H}[n] \) as an iterated blowup of a mixed Tate variety along a mixed Tate locus and the blowup formula for motives. The varieties \( \mathcal{X}_{H}[n] \) are defined over \( \mathbb{Z} \), so we only need to worry about the algebraic differential form \( \phi_{+}(\eta_{\Gamma}) \), which may be defined over \( \mathbb{C} \) or possibly a smaller field.

**Remark 6.7.** This renormalization method does not provide a direct description of the period integral \((6.6)\) as a combination of multiple zeta values with coefficients in \( \mathbb{Q}((2\pi i)^{-1}) \). In fact, one knows by \([3]\) that periods of mixed Tate motives over \( \mathbb{Z} \) would be of this form and that the varieties \( \mathcal{X}_{H}[n] \) are defined over \( \mathbb{Z} \), but we have constructed the differential form \( \eta_{\Gamma} \) as a form cohomologous to \( \pi_{n}^{*}(\omega_{\Gamma,\mathcal{C}}) \) in \( H^{*}(\mathcal{X}_{H}[n], D_{n}, \mathbb{C}) \). Thus, a priori, \( \eta_{\Gamma} \) and its renormalization \( \phi_{+}(\eta_{\Gamma}) \) are only \( \mathbb{C} \)-forms and not \( \mathbb{Q} \)-forms. A more detailed analysis of the relation to multiple zeta values will be carried out elsewhere \([6]\).
6.3.4. The massive case. In this setting of regularization by algebraic forms with logarithmic poles on \( X_H[P] \) and renormalization by pole subtraction and Birkhoff factorization of the regularized forms, one can treat the massive case as a formal series, using an expansion such as (4.18)

\[
\omega_{\Gamma,m,C} = \sum_{\ell} \prod_{e \in E_\Gamma} \frac{B_{\Gamma,m,D,\ell}^C}{(2\pi)^D} \beta_{\Gamma,m,D,\ell}^C
\]

or the asymptotic expansion as in (4.27)

\[
\omega_{\Gamma,m,C} \sim \sum_{\ell} C_{\Gamma,m,D,\ell} \omega_{\Gamma,m,D,\ell} C,
\]

One can then consider a regularization \( \eta_{\Gamma,m,D,\ell} \) of the forms \( \beta_{\Gamma,m,D,\ell}^C \), or, respectively, of the \( \omega_{\Gamma,m,D,\ell} C \) and then proceed to the pole subtraction and renormalization of the forms \( \eta_{\Gamma,m,D,\ell} \) according to the Rota–Baxter renormalization procedure described in this section.

6.4. The renormalization group. In the Connes–Kreimer theory of renormalization based on Birkhoff factorization [9], the beta function of the renormalization group equation is lifted to the Lie algebra of the affine group scheme dual to the Hopf algebra of Feynman graphs. In the subsequent work [10], it is further lifted to the level of a universal Hopf algebra, where the universal counterterms are obtained from a universal singular frame depending on the beta function, seen as an element in a free Lie algebra. This algebraic approach to the renormalization group flow was further developed in terms of Dynkin idempotents, [12], see also [21] for a survey of these recent results.

Let \( A \) be an associative algebra. The Dynkin operator \( D_n \in \mathbb{Q}[S_n] \) with \( S_n \) the symmetric group, is defined by the property that

\[
[x_1, x_2, \ldots, x_n] = x_1 x_2 \cdots x_n D_n,
\]

for any collection of elements \( x_1, \ldots, x_n \in A \), with the symmetric group acting on the right. One defines \( D = \sum_{n \geq 0} D_n \).

Consider then the group of algebraic Feynman rules \( G(\mathcal{R}) = \text{Hom}_{\text{Alg}}(\delta, \mathcal{R}) \) and let \( g(\mathcal{R}) \) be the corresponding Lie algebra. It is shown in [12] that right composition with the Dynkin operator \( D \) gives a bijective map from \( G(\mathcal{R}) \) to \( g(\mathcal{R}) \), whose inverse map is given by the “universal frame for renormalization” in the sense of [10],

\[
g(\mathcal{R}) \ni \beta \mapsto \sum_{n \geq 0, k_i > 0} \frac{\beta_{k_1} \cdots \beta_{k_n}}{k_1(k_1+k_2) \cdots (k_1+\cdots+k_n)} \in G(\mathcal{R}),
\]

where \( \beta = \sum_n \beta_n \) with respect to the grading of the Hopf algebra.

In our setting we obtain an interpretation of the renormalization group by applying this construction with \( \mathcal{R} \) the Rota–Baxter algebra of differential forms on the wonderful configuration spaces with logarithmic poles along their boundary divisors. Namely, one identifies the beta function of the renormalization group with that element \( \beta \) in the Lie algebra \( g(\mathcal{R}) \) that is the image under right composition with the Dynkin operator \( D \) of the negative piece of the Birkhoff factorization

\[
\phi_-(\Gamma) = -T(\eta_{\Gamma} + \sum_{\gamma \subset \Gamma} \phi_-(\gamma) \wedge \eta_{\Gamma/\gamma})
\]

of the regularized Feynman amplitude \( \eta_{\Gamma} \). Thus, \( \beta \in g(\mathcal{R}) \) is determined by the property that

\[
\phi_- \cdot D = \beta \quad \text{or equivalently} \quad \phi_- = \sum_{n \geq 0, k_i > 0} \frac{\beta_{k_1} \cdots \beta_{k_n}}{k_1(k_1+k_2) \cdots (k_1+\cdots+k_n)}.
\]
7. Beyond Scalar QFTs

The more realistic QFTs, such as gauge theories, involve particles with spin or polarization and the propagators for such fields are slightly different from the scalar case. We show that, nonetheless, these propagators have similar properties, which are sufficient to describe the corresponding Feynman integrals, as in scalar case, as periods.

7.1. Spin 1/2 particles: the Dirac fermions. In Euclidean quantum electrodynamics, the propagator for a Dirac field representing the electron is the fundamental solution of the Dirac equation

\[ (-i\not\partial + m)S(x) = \delta(x) \]

where \( \gamma^\mu \) are the gamma matrices and \( \not\partial = \gamma^\mu \partial_\mu \). The Green function \( S(x) \) can be given as

\[ S_{m,D}(x) = \frac{1}{(2\pi)^D}\int_{\mathbb{R}^D} dk \frac{e^{ik\cdot x}}{-k + m} = \frac{1}{(2\pi)^D}\int_{\mathbb{R}^D} dk \frac{e^{ik\cdot x} \sqrt{k + m}}{\|k\|^2 + m}. \]

**Lemma 7.1.** The Euclidean Dirac propagator on \( \mathbb{R}^D \), with \( D = 2\lambda + 2 \), satisfies

\[ S_{m,D}(x) = \left( \frac{1}{2\pi} \right)^{\lambda+1}(i\gamma^\mu x_\mu)\frac{m^2}{\|x\|^{\lambda+1}}\left( \frac{\lambda}{\|x\|}K_\lambda(\sqrt{m}\|x\|) + \frac{\sqrt{m}}{2}(K_{\lambda-1}(\sqrt{m}\|x\|) + K_{\lambda+1}(\sqrt{m}\|x\|)) \right) \]

\[ + \left( \frac{1}{2\pi} \right)^{\lambda+1}\frac{m^{\lambda+2}}{\|x\|^\lambda}K_\lambda(\sqrt{m}\|x\|). \]

**Proof.** A simple computation using (7.1) gives the Dirac propagator \( S_{m,D} \) on \( \mathbb{R}^D \) in terms of the scalar field propagator as follows:

\[ S_{m,D} = (-i\not\partial + m)G_{\sqrt{m},D}. \]

Then, we substitute the scalar propagator \( G_{m,D} \) given in Lemma 4.6 and we use the identity

\[ \frac{\partial}{\partial z}K_\lambda(z) = -\frac{1}{2}(K_{\lambda+1}(z) + K_{\lambda-1}(z)). \]

**Corollary 7.2.** The Dirac propagator \( S_{m,D}(x) \) is smooth on \( \mathbb{R}^D \setminus \{0\} \) with \( D = 2\lambda + 2 \). Moreover, it diverges at \( x = 0 \), and \( S_{m,D}(x) \to 0 \) as \( \|x\| \to \infty \).

7.2. Spin 1 particles: the gauge bosons. The propagator of a spin 1 boson with mass \( m \) is given by

\[ \Delta^\mu_\nu(x) = \frac{1}{(2\pi)^D}\int_{\mathbb{R}^D} dk \frac{e^{ik\cdot x}}{\|k\|^2 + m} \left( \frac{g^\mu_\nu + k^\mu k^\nu}{m^2} - \frac{k^\mu k^\nu}{m^2 + m^2} \right) \]

in the Stueckelberg gauge, depending on a gauge parameter \( \alpha \).

**Lemma 7.3.** The propagator for a massive vector field can be derived from the scalar field propagator \( G_{m,D} \):

\[ \Delta^\mu_\nu_{m,D}(x) = g^\mu_\nu G_{\sqrt{m},D}(x) + \frac{1}{m^2} \left( \partial^\mu \partial^\nu G_{\sqrt{m},D}(x) - \partial^\mu \partial^\nu G_{\sqrt{m},D}(x) \right). \]

Therefore, the propagator \( \Delta^\mu_\nu_{m,D}(x) \) is smooth on \( \mathbb{R}^D \setminus \{0\} \) with \( D = 2\lambda + 2 \). Moreover, it diverges at \( x = 0 \), and \( \Delta^\mu_\nu_{m,D}(x) \to 0 \) as \( \|x\| \to \infty \).

In other words, the Dirac propagators and the massive gauge bosons propagators can be related to propagators for the scalar case, hence any Feynman integrand involving an amplitude constructed out of these two types of propagators can also be interpreted in terms of periods of mixed Tate motives, following the arguments we developed for the scalar case.
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Faculté des Sciences, de la Technologie et de la Communication, University of Luxembourg, 6 rue Richard Coudenhove-Kalergi, L-1359 Luxembourg

E-mail address: ozgur.ceyhan@uni.lu

Division of Physics, Mathematics and Astronomy, Mail Code 253-37, Caltech, 1200 E. California Blvd. Pasadena, CA 91125, USA

E-mail address: matilde@caltech.edu