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Resummation for quantum propagators in bounded spaces

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I. INTRODUCTION

The quantum mechanical propagator $K(y, x; T)$ (also called the kernel) is defined in terms of the configuration space matrix elements of the time evolution operator,

$$K(y, x; T) = \langle y|e^{-iH T}|x\rangle \theta(T),$$

(1)

where $H = \frac{p^2}{2m} + V(x)$ is the Hamiltonian of a system with potential $V$. This kernel is a Green function for the position space Schrödinger equation with the boundary condition $K(y, x; 0^+) = \delta^D(y - x)$ and as such contains the full information about the quantum system. Except in some specific cases, or for systems with special symmetries, finding the explicit form of the kernel can be very nontrivial and even conventional perturbation theory has its limitations, especially for strongly perturbing potentials, so new techniques for its calculation or estimation are of significant interest.

In [1] two invertible integral transforms of the kernel were introduced. The “hit function” defined there is denoted by $\tilde{H}(z|y, x; T)$ and represents the contribution to the kernel involving propagation that passes through (“hits”) the spatial point $z$; it can be used to recover the propagator by integrating over this point:

$$\tilde{H}(z|y, x; T) = \frac{1}{T} \int_0^T dt K(z, x; t)K(y, z; T - t),$$

(2)

$$K(y, x; T) = \int d^Dz \tilde{H}(z|y, x; T).$$

(3)

This function also finds application in prior studies of Brownian motion and Levy random walks where it is related to the concept of “local time” [2,3] that measures the proportion of time for which a particle’s trajectory is found in a particular region. For the free particle, which will be the focus of this paper, the hit function was calculated in one dimension in the imaginary time formalism (“Euclidean hit function”) in [1] to be

$$\tilde{H}_0(z|y, x; T) = \frac{m}{2T} \text{erfc} \left[ \frac{m}{2T}(|x - z| + |z - y|) \right], \quad T \geq 0,$$

(4)
which was verified by a numerical sampling of the Feynman path integral (note that this function is constant for \( z \) between \( y \) and \( x \), and erf is the complementary error function defined precisely below).

For potentials that are sharply localized, such as\(^1\) \( V(x) \sim \lambda \delta_\epsilon(x-a) \), the perturbative expansion of the path integral representation of the kernel (see Eq. (6) below) in powers of the coupling constant, \( \lambda \), naturally constrains the path integral to trajectories that must pass through (or in a neighborhood of) the point \( a \) various times. This is the case, for example, for many scattering problems, delta function potentials as studied by Grosche [4–6] and contact interactions [7–9], and varied models of lattice structure in condensed matter [10,11]. In the relativistic case, constrained path integrals arise in the context of the Casimir effect where one is interested in trajectories that touch the conducting plates [12–15], and for Dirichlet boundary conditions imposed in some spatial region [16–25] where the contributions from paths passing through this region should be removed. A fresh approach to calculating, either analytically, approximately or numerically, such path integrals will therefore find wide application.

The purpose of this article is threefold. We will introduce the generalization of the hit function to the “\( n \)-hit function” that counts trajectories that pass through \( n \) intermediate spatial points \( \{z_1, \ldots, z_n\} \) in chronological order whilst traveling between the endpoints \( x \) and \( y \). This \( n \)-hit function can be defined as the constrained path integral (again, in Euclidean space, with \( \tau_{n+1} := T \))

\[
\mathcal{H}(z_1, \ldots, z_n | x, y; T) := \int_{\mathcal{D}(0)=\emptyset} \mathcal{D}\mathcal{X}(\tau) \prod_{i=1}^{n} \int_{0}^{\tau_i} d\tau_i \delta^{D}[x(\tau_i) - z_i] \times e^{-\int_{0}^{T} d\tau \left( \frac{\partial^{2}}{2} + V(x(\tau)) \right)} \theta(\tau). \tag{5}
\]

Unless otherwise indicated we will mean the free case when we refer to the \( n \)-hit function which is made explicit in expressions via the notation \( \mathcal{H}_0 \). We will uncover relations between the \( n \)-hit function for different order, \( n \), and in different number of spatial dimensions, \( D \), and will give explicit formulas that allow it to be determined for arbitrary \( n \) and \( D \). We will also use this function and the theory of Padé approximants to conjecture a closed analytical formula for the propagator in the presence of Dirichlet boundary conditions in arbitrary geometries and support our claim by applying it to some simple examples. In practical calculations, this reflecting is approached by a sequence of approximants, which essentially represent a perturbative expansion of the kernel appropriate for the strong coupling limit.

It is important to highlight some prior studies based on a similar philosophy. The path decomposition expansion of [26–29] splits up the path integral representation of the Green function associated to the kernel in terms of contributions from trajectories in different spatial regions; in [30,31] the multistep propagator is introduced with the aim of eliminating a possibly dangerous overcounting of certain types of trajectory in the path integral; the effective action of a (relativistic) scalar field in a localized potential is found in [32] by resumming a perturbative expansion of the interaction (in fact our \( n \)-hit function is related to an infinite set of intermediate functions that enter this resummation); and finally, Polyakov long ago formulated a representation of the (one-particle-reducible) contributions to the \( N \)-point configuration-space amplitudes in \( \phi^4 \) theory in terms of a relativistic point-particle path integral constrained to pass through the prescribed external points [33].

However, here we work with the path integral representation of the propagator as is, in configuration space, and work out the relative contributions from the (infinite number of) trajectories that pass through the prescribed points. The \( n \)-hit function we define has a simple physical interpretation according to context: in the quantum mechanical sense it represents the (un-normalized) amplitude of propagation for a particle to travel from point \( x \) to point \( y \) in time \( T \), going through the points \( z_1, \ldots, z_n \) (in time order) per unit \((n+1)\)-volume \( d^{D}z_1 \cdot \ldots \cdot d^{D}z_n d^{D}y \) (following the usage in [34] we will refer to such an amplitude as a relative amplitude); in the context of Brownian dynamics it describes the probability density of the propagation taking place under the same constraints, per unit \( n \)-volume of the intermediate points.

This information can be useful for analytic and numerical estimations of the kernel for systems whose path integral cannot be computed in closed form, since it can indicate which trajectories, or regions of space, will be dominant in its determination. In fact the hit function was already sampled in works based on worldline numerics [1] and the extension considered here could shed light on the undersampling problem encountered in [35–37]. Looking further ahead, the methods we describe here are well adapted to the problem of estimating the Casimir energy for arbitrary surface geometries, even for partially conducting plates. Indeed, via the lesser well-known worldline formalism of quantum field theory [38–40], the generalized hit function defined here admits a clear extension to field theory processes where it will be defined for relativistic point particles in Minkowski space (as an immediate example, the \( n \)-hit function can be converted into Polyakov’s representation of \( \phi^4 \) amplitudes cited above by transforming to Minkowski space and introducing a (Schwinger) integral over the trajectories’ proper time [33]). This case shall be addressed in future work, but see [22] where the heat kernel with Dirichlet boundary conditions on a \( D \)-dimensional manifold (\( D \)-ball) was obtained using worldline techniques and conformal transformations.

The outline of this paper is as follows: Sec. II introduces the \( n \)-hit function in terms of path integrals and gives also an \( n \)-fold integral representation in terms of the free particle kernel \( K_0 \). Section III presents one of the main results of this work, a conjectured representation for the propagator or heat kernel on a bounded region with Dirichlet boundary conditions based on Padé approximants. Section IV A shows a closed formula for the free \( n \)-hit function in \( D = 3 \) spatial dimensions, the case of most direct physical relevance and also special for being easily accessible analytically. Section IV B gives the energy representation of the general \( n \)-hit function for any space dimension \( D \). Section VI C presents an

\(^1\)Here \( \delta \) indicates a localized, \( \delta \)-function-like potential that may be smeared over some small characteristic scale \( \epsilon \) (it may be extended in a codimension \( d > 0 \) subspace of \( \mathbb{R}^d \)).
analytical example where a proof of the conjectures presented in Sec. III is given in a simple one-dimensional example. Section V demonstrates how certain general operations connect $n$-hit functions in different spatial dimensions $D$ and of different orders $n$ to one another. An explicit realization of these operators is given, providing a set of identities relating the functions that in turn leads to a hierarchy which generates all possible $n$-hit functions. Section VI A contains numerical results for the propagators inside a sphere in $D = 3$ and an infinite plane (also in $D = 3$) using the Padé representation proposed in Sec. III. In Sec. VI B we obtain an integral energy-representation of the general $n$-hit function in any $D$. Section VII presents our conclusions and insight on further applications of the method proposed and more generally of the $n$-hit function and its connection to Feynman’s description of the path integral. Finally, for the sake of clarity most detailed calculations and prior results have been deferred to Appendices.

II. QUANTUM PROPAGATION AND THE $n$-HIT FUNCTION

The problem we will address is that of quantum propagation in the presence of a potential, and while the methods outlined here can be applied to more general situations we will mainly focus on a special kind of singular potential with support on a predetermined surface, $V(x) = \lambda \int_D \delta(x-z) d\sigma$. Here $z$ is a position vector on $S$, $d\sigma$ represents a measure that fully parametrizes the surface $S$ and $\lambda$ is a positive constant that fixes the strength of the potential. The propagation amplitude of a nonrelativistic particle between position $x$ at time $t = 0$ and position $y$ at time $t = T$ in the presence of a potential is known to be given in natural units by the path integral

$$ (y; T|x; 0) = \langle y|e^{-i\hat{H}T}|x \rangle = \int_{\mathcal{X}(0)=x}^{\mathcal{X}(T)=y} \mathcal{D}X(T) e^{-\int_0^T dt \{ i\mathcal{H}[X(t)] \}}. $$

(6)

Henceforth we will Wick rotate to the imaginary time formulation of quantum mechanics in which quantum free propagation becomes Brownian motion and all amplitudes become real. This also has the effect of ensuring that the hit functions we are to define will be real valued. It is possible to recover the quantum mechanical propagation amplitudes by reversing the Wick rotation afterwards, and as our results will show the $n$-hit functions can be analytically continued back to the real time axis. However, we point out that proper care must be taken due to the presence of branch cuts in the complex $T$-plane and draw attention to the fact that in the real time formalism the generalized hit functions are complex-valued. Motivated by future applications we will usually henceforth set $m = 1/2$ – if desired the mass $m$ can always be recovered by dimensional analysis in all our results, (the choice $m = 1/2$ corresponds to a unit diffusion coefficient) which converts (1) to

$$ K(y, x; T) = \int_{\mathcal{X}(0)=x}^{\mathcal{X}(T)=y} \mathcal{D}X(T) e^{-\int_0^T dt \{ i\mathcal{H}[X(t)] + V(x(t)) \} \theta(T)}. $$

(7)

A standard way to deal with the localized potential considered here is to expand the exponential to obtain a perturbative series in the constant $\lambda$,

$$ K(y, x; T) = \int_{\mathcal{X}(0)=x}^{\mathcal{X}(T)=y} \mathcal{D}X(T) e^{-\int_0^T dt \{ \lambda \int D^Dz \delta(x-z) \} \theta(T)} \times \left( \sum_{n=0}^{\infty} \frac{(-1)^n}{n!} \int_0^T V(x(t)) d\tau \right)^n \theta(T). $$

(8)

Hence the $n$th term of this series requires the determination of a constrained path integral

$$ \mathcal{I}_n([z_i], \{ \tau_i \}) := \int_{\mathcal{X}(0)=x}^{\mathcal{X}(T)=y} \mathcal{D}X(T) e^{-\int_0^T dt \{ \lambda \int D^Dz \delta(x-z) \} \theta(T)} \prod_{j=1}^{n} \delta^D(x(\tau_j) - z_i). $$

(9)

Remarkably, following the steps in [1] one can determine this in terms of the free particle propagator (see also Appendix B), $K_0$, as

$$ \mathcal{I}_n([z_i], \{ \tau_i \}) = K_0(x, z_{\tau(1)}; \tau_{\tau(1)}) \cdot K_0(z_{\tau(1)}, z_{\tau(2)}; \tau_{\tau(2)} - \tau_{\tau(1)}) \cdots K_0(z_{\tau(n-1)}, y; T - \tau_{\tau(n-1)}). $$

(10)

Here $\Pi$ is the time-ordering permutation of subindices defined so that $0 \leq \tau_{\tau(1)} \leq \tau_{\tau(2)} \leq \cdots \leq \tau_{\tau(n)} \leq T$, and in $D$ spatial dimensions the free propagator is

$$ K_0(x, y; \tau) = \frac{1}{(4\pi \tau)^{D/2}} e^{-(x-y)^2/4\tau} \theta(\tau). $$

(11)

The physical meaning of $\mathcal{I}_n$ is manifest: it represents the relative amplitude\(^3\) of propagation starting off from $x$ at time $\tau = 0$ and arriving at $y$ at time $\tau = T$ going through the prescribed intermediary points $z_1, z_2, \ldots, z_n$ at their assigned times $\tau_1, \tau_2, \ldots, \tau_n$ in time-ordered fashion.

For the potential $V(x)$ with support on $S$ the full series (8) can be found from (10)

$$ K(y, x; T) = \sum_{n=0}^{\infty} \int_S d\sigma_n \cdots \int_S d\sigma_1 \frac{(-\lambda)^n}{n!} \int_0^T d\tau_n \cdots \int_0^T d\tau_1 \times \prod_{\pi} \mathcal{I}_n([z_{\tau(\pi)}], \{ \tau_i \}) \theta(T). $$

(12)

where $\pi \in S_n$ denotes a permutation in the symmetric group $S_n$ and $\sum_\pi$ denotes sum over all such permutations. Notice that the same method can be used for a more general potential by rewriting $V(x) = \int V(z) \delta(x-z) d^Dz$ and doing the appropriate replacements, we discuss this matter in Sec. III B. This is all well known and constitutes one approach to arriving at

3A relative amplitude $\Psi$ is an amplitude density per unit Vol\(^{+1}\) normalized so that $\int d^Dz_1 \cdots d^Dz_n d^D\Psi = 1$. We refer to these contributions as amplitudes even though they are actually the Euclidean version (probabilities) of a true quantum mechanical amplitude.
standard perturbation theory in quantum mechanics [34], and it is at this point that we will introduce a series of developments aimed at treating this problem in a distinct manner.

First, we point out that the integrations in (12) contain all possible time orderings of the variables \(\tau_i\), but due to permutation symmetry each ordering yields the same contribution after integrating over \(S\). This allows us to fix an arbitrary ordering and cancel the \(1/n!\) prefactor as we do below in (15). For this reason it is now convenient to choose a canonical ordering \(\tau_0 \leq \tau_1 \leq \tau_2 \leq \cdots \leq \tau_n \leq \tau_n+1 := T\) and define, in accordance with (5), the function

\[
H_0(z_1, z_2, \ldots, z_n; y, x; T) := \int_0^T d\tau_n \int_0^{\tau_n} d\tau_{n-1} \cdots \int_0^{\tau_1} d\tau_1 \ K_0(x, z_1; \tau_1) \\
\times K_0(z_1, z_2; \tau_2 - \tau_1) \cdots K_0(z_n, y; \tau_n+1 - \tau_n),
\]

that we shall call the \(n\)-hit function: it relates to the origin of the propagator (2) through \(H_0(z; y, x; T) = T^{-1/2}H_0(z; y, x; T)\), which illustrates that in this article we have chosen a different normalization convention than in [1]. For notational convenience we also drop the bar used there to distinguish between normalized and un-normalized probability distributions. The path integral representation of this function is nothing other than (5) with the potential \(V(x) = 0\) [indicated in (13) by the subscript 0].

The \(n\)-hit function can be normalized into a relative amplitude\(^4\) by noticing that

\[
\frac{n!}{T^n} \int_{\mathbb{R}^D} \cdots \int_{\mathbb{R}^D} H_0(z_1, \ldots, z_n; y, x; T) = K_0(y, x; T),
\]

which also provides the inverse integral transform that recovers the original propagator from the \(n\)-hit function. In this manner it is consistent to take the 0-hit function to be the free space propagator, \(K_0(y, x; T)\).

Through direct substitution we can now write the propagator (12) in a perturbative series that involves successive hit functions integrated over \(S\):

\[
K(y, x; T) = K_0(y, x; T) - \lambda \int_S d\sigma_1 \ H_0(z_1; y, x; T) \\
+ \lambda^2 \int_S d\sigma_1 d\sigma_2 \ H_0(z_1, z_2; y, x; T) + \cdots. \tag{15}
\]

This power series in \(\lambda\) constitutes the basis of our approach to calculate the propagator and corresponds to counting contributions scattered by the potential exactly once, twice, and so forth. This representation makes clear that knowledge of just the free hit functions is sufficient to recover the kernel in the presence of a localized potential.

In dimension \(D\) odd, we will show that it is possible to evaluate \(H_0\) in closed form for an arbitrary number of intermediate points \(n\). For the case of even \(D\), although it is straightforward to calculate the hit function using the various relations we will outline below, a closed formula of similar simplicity has proven elusive so far, except for the special cases implied by (67). These and other results regarding the \(n\)-hit function will be postponed to a later section, since instead we will focus first on the case \(D = 3\) of special physical significance.

III. PADÉ REPRESENTATION OF THE PROPAGATOR

The problem of calculating the propagator in quantum mechanics, or equivalently the heat kernel, in a bounded region of Euclidean space and on compact manifolds has been investigated for a long time and consequently an arsenal of methods and remarkable results are already known [41]. With the exception of systems with a certain degree of symmetry the propagator is not known in analytic form and therefore approximations or numerical calculations must be made. A solution of the Dirichlet problem in terms of analytic expressions will be relevant in light of the importance of heat kernel methods used in quantum field theory [19,21,22]. In this section we propose an analytical solution and claim that this introduces a representation for the propagator in bounded regions, or equivalently an analytical expression for the path integral computed from trajectories that do not intersect a given boundary. We lend support to this claim with numerical evidence in later sections.

A. Dirichlet boundary conditions

We are interested in the solution of the Dirichlet problem, that is, the calculation of the propagator subject to the amplitudes vanishing on the boundary \(S\). This condition can be imposed in the propagator by taking the limit \(\lambda \to +\infty\) that makes the potential barrier impenetrable. It is therefore necessary to determine the limiting behavior as \(\lambda \to +\infty\) of the perturbative series representation of the propagator in (15), which naturally falls well outside the radius of convergence of the series. Grosche has demonstrated that in the presence of a Dirac delta potential in \(N\) points the energy representation of the quantum mechanical propagator is exactly a rational function of the coupling strength \(\lambda\) in dimensions \(D = 1, 2, 3\) [4–6]. We propose that while the exact propagator will not (usually) be a rational function for a delta potential with support on a surface \(S\), the propagator can nevertheless be approximated by a sequence of rational functions of \(\lambda\) to any desired degree of precision. Following this notion we will turn to Padé approximants and use an idea introduced by Bender and Boettcher in [42] whereby the value of a function at infinity can be obtained under favorable circumstances from the coefficients of its perturbative series using diagonal Padé approximants.

As a result, our proposed method is to determine a sequence of diagonal Padé approximants, \(P_N^n\) (see Appendix E), to the series representation of the propagator based on the iterated integrals of the hit function and then to extrapolate
to the strong coupling limit, or \( N \to \infty \). As we show in the Appendix, these can be expressed as a quotient of determinants that give a rational approximation to the propagator. Following this approach, after calculating the Padé approximant to \( (15) \) we are led to conjecture that the difference between the propagator and the propagator in free space is given by the limit of a quotient of determinants

\[
K(y, x; T) - K_0(y, x; T) = \lim_{N \to \infty} \frac{\prod_{n=0}^{N} c_n}{\prod_{n=0}^{N} c_n},
\]

where

\[
c_n := (-1)^{n+1} \int_S d\sigma_1 \cdots \int_S d\sigma_{n+1} H_0(z_1, \ldots, z_{n+1} | y, x; T).
\]

While we have found strong support for this claim in our numerical calculations, we do note that it remains a conjecture at this point. However it is important to emphasize the significance of \( (16) \) and \( (17) \) as they state how the effect of the boundary on the propagator can be accounted for purely in terms of multiple scatterings at the boundary and that it is calculated explicitly by integrating hit functions over that boundary. To the best of the authors’ knowledge this is a novel development.

Let us also draw attention to the fact that, although completely different in content, \( (16) \) resembles in spirit the original formulation of the Gelfand-Yaglom theorem [43] for a Schrödinger problem with Dirichlet boundary conditions in which the effect of a potential on the path integral normalization is expressed as a finite quotient of functional determinants that were calculated through a limiting process from finite dimensional matrix determinants.

Together \( (16) \) and \( (17) \) state one of the principal results of the present work, while simultaneously motivating us to determine the analytic form of the free \( n \)-hit functions for their practical use in \( (16) \). It is also clear from a practical point of view that a finite \( N \) must be chosen at some point to carry out an actual approximate calculation of the propagator. In this sense our result also constitutes a pragmatic approximation that in contrast to conventional series representations should get more accurate the stronger the coupling.

### B. Other static potentials

Before continuing, we pause to note that an arbitrary potential \( V(x) \) may not be tractable with the present approach as the sequence of Padé approximants may not converge, yet we wish to illustrate that a wide family of potentials could still be treated with this method. We start by writing the series expansion \( (8) \) for the potential

\[
V_L(x) = \lambda \int d^D z \delta(x-z).
\]

Following the procedure outlined is tantamount to the formal substitution \( \int_S d\sigma_l \to \int d^D z \delta(z_i) \) everywhere. In this case the parameter \( \lambda \) could be an appropriate coupling extracted from the physical potential or could be artificially introduced; in the latter case after calculating the perturbative series analog to \( (15) \) the limit \( \lambda \to 1 \) is taken. The following diagonal Padé representation for the propagator is obtained:

\[
K(y, x; T) = \lim_{N \to \infty} \frac{\prod_{n=0}^{N} c_n}{\prod_{n=0}^{N} c_n},
\]

where \( \varphi = \sum_{p=0}^{L} C_p \) and

\[
K_0 = K_0(y, x; T),
\]

\[
K_n = (-1)^n \int d^D z_1 \cdots \int d^D z_n V(z_1) \cdots V(z_n) \times H_0(z_1, \ldots, z_n | y, x; T), \quad n \geq 1.
\]
In contrast to (16) we may now also consider the nondiagonal Padé approximants: From the numerical standpoint this is often useful to speed up convergence as they provide more information on the convergence behavior of the sequence, in which case the analog to (19) can be easily obtained from (E5) by taking $M \neq N$. It will be interesting to investigate the kind of potentials that this method can be applied to. Although the results presented below lend weight to the argument that this approach will be well suited to localized potentials, the authors venture that it in principle may be applicable as a kind of potentials that this method can be applied to. Although

In the general case the evaluation of the hit function defined in (13) is greatly facilitated by using the “energy representation” of the kernel, its Fourier transform with respect to transition time $T$,

$$
\hat{K}(y, x; \omega) = \int_{-\infty}^{\infty} dT \, K(y, x; T) e^{i\omega T}.
$$

Here we use this representation of the kernel to manipulate the integral representation of the $n$-hit function into the form of an inverse Mellin transform.

We begin with the observation that the $n$-hit function can be written as (in fact, this result holds for the hit function defined with respect to an arbitrary background potential by replacing $K_0 \rightarrow K_0$; cf. footnote in Sec. II)

$$
H_0(z_1, \ldots, z_n; y, x; T) = \prod_{i=1}^{n} \int d^2 z_i \, H_0(z_1, \ldots, z_i; y, x; T) = 1
$$

[cf. (14)].

---

It is pleasant to note that this property is the time analogue of the reproducing kernel in space for the same function:

$$
\int_{\mathbb{R}^2} f_{y-i}(T - \tau) f_{x-i}(\tau) \, d^2 \tau = \pi^{3/2} f_{y-x}(T).
$$

Here space vector spaces are added; over time it is their lengths that are to be added. Both behaviors can be traced back to the fact that free propagation forgets about the prior points it has previously passed through and is isotropic so that the amplitude depends not on the relative orientations of the edges of the path between the points $z_i$, but rather on their lengths only.

---

IV. DETERMINING THE $n$-HIT FUNCTION

Based on the proposal outlined above, in this section we calculate the $n$-hit function in the special case of dimension $D = 3$ and then give a general formula, (29), in arbitrary dimension. From these results we derive relations between the $n$-hit functions for differing $D$ and $n$ which allow us to transform between different values of these parameters with simple operations on the functions.

A. Three-dimensional hit function

It turns out that $D = 3$ dimensions is a special case for which the iterated integrals in (13) can be calculated straightforwardly. Using the relations in Appendix A, the key observation is that the function $f_0(T) = \frac{1}{1 - e^{-\alpha T}}$ is a reproducing kernel over $T$ in the sense that

$$
\int_0^{T} f_0(T - \tau) f_0(\tau) \, d\tau = \frac{\alpha}{\alpha + \beta} f_{\alpha+\beta}(\tau),
$$

so that in terms of the (finite-interval) convolution $[f_\alpha \ast f_\beta](T) = \sqrt{\frac{4\pi}{\alpha \beta}} f_{\alpha+\beta}(\tau)$. This result allows us to perform all integrals recursively in (13) to determine the $n$-hit function for arbitrary order.

Repeated application of the convolution property and the results for the integrals given in Appendix A yield a closed form for the $n$-hit function in $D = 3$:

$$
H_0(z_1, \ldots, z_n; y, x; T) = \frac{1}{(4\pi)^{\frac{n+3}{2}}} \frac{\Delta}{\prod_{\alpha=1}^{n+1} \Delta_\alpha} \theta(T),
$$

where we introduced $\Delta_k := |z_k - z_{k-1}|$ and defined $z_0 := x$ and $z_{n+1} := y$, so that $\Delta = \sum_{i=1}^{n+1} \Delta_i$ is the total length of the polygonal path from $x$ to $y$ with prescribed vertices $\{z_i\}$. We shall make frequent use of this notation for the remainder of this article.

To illustrate this result we give the explicit form of the first two hit functions as

$$
H_0(z, y; x; T) = \frac{1}{(4\pi)^{3/2}} \frac{e^{-\frac{\Delta}{4\pi^{3/2}}}}{\Delta_1 \Delta_2} \theta(T),
$$

and

$$
H_0(z_1, z_2; y, x; T) = \frac{1}{(4\pi)^{5/2}} \frac{e^{-\frac{\Delta}{4\pi^{5/2}}}}{\Delta_1 \Delta_2 \Delta_3} \theta(T).
$$

B. Arbitrary $n$-hit function in dimension $D$

In the general case the evaluation of the hit function defined in (13) is greatly facilitated by using the “energy representation” of the kernel, its Fourier transform with respect to transition time $T$,

$$
\hat{K}(y, x; \omega) = \int_{-\infty}^{\infty} dT \, K(y, x; T) e^{i\omega T}.
$$

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the $t_i$ now represent the lengths of proper-time intervals. Using now the well-known Fourier representation of the Dirac delta distribution $\delta(t) = \frac{1}{\pi} \int_{-\infty}^{\infty} d\omega \ e^{i\omega t}$ the convolution theorem yields the integral representation (again, replacing the free propagators on the right-hand side by their counterparts corresponding to some other background potential yields an analogous integral representation of the $n$-hit function for that potential)

$$H_0(z_1, z_2, \ldots, z_n[y, x; T]) = \int_{-\infty}^{\infty} \frac{d\omega}{2\pi} e^{-i\omega T} \tilde{K}_0(y, z_n; \omega) \cdots \tilde{K}_0(z_2, z_1; \omega) \tilde{K}_0(z_1, x; \omega).$$

(26)

A physical interpretation of (26) is that as the particle is scattered multiple times in vacuum (or in a static potential) the amplitudes corresponding to a given fixed energy are multiplied, that is, the scatterings happen independently and do not change the energy; we then calculate the total amplitude by integrating over the values of energy at which this scattering took place. In our conventions, (26) contains the Fourier transform of the free propagator according to (24) which for the free particle is known, for dimension $D$, in terms of the Macdonald function (modified Bessel function of the second kind), $K_{D/2}$ (see [45], for example). Indeed, with $\nu := \frac{2-D}{2}$ one obtains for arbitrary $D$,

$$\tilde{K}_0(y, x; \omega) = \frac{2}{(4\pi)^{D/2}} \left(\frac{|x-y|}{2\sqrt{-i\omega}}\right)^\nu K_{\nu}(\sqrt{-i\omega}|x-y|).$$

(27)

Substituting this into (26) into we arrive at the following Fourier integral representation:

$$H_0(z_1, \ldots, z_n[y, x; T]) = C_{D,n} \int_{-\infty}^{\infty} \frac{d\omega}{2\pi} \left(e^{-i\omega T}\right) K_{\nu}(\sqrt{-i\omega}|x-z_1|) \cdots K_{\nu}(\sqrt{-i\omega}|z_n-y|),$$

(28)

where we introduced the prefactor

$$C_{D,n} = (2\pi)^{-D(n+1)/2} (|x-z_1| \cdots |z_n-y|)^\nu.$$

We now introduce the complex variable $\zeta = \sqrt{-i\omega}$ and proceed to deform the contour of integration using Cauchy’s theorem; when $T \ll 0$ we can collapse the path onto the real axis on the complex $\zeta$ plane, leading to the vanishing of the integral, as required by the Heaviside $\theta$ functions in (11), whereas in the case $T > 0$ the contour can no longer be deformed in this manner but rather can be deformed to run parallel to the imaginary axis shifted a small amount $\theta^+$ towards the right in accordance with the Feynman prescription;

7The relevant identity is (see, for example, [44])

$$\int_0^\infty dt e^{-\beta t} = 2\left(\frac{\beta}{\gamma}\right)^{\nu/2} e^{-\beta|\gamma|} K_{\nu}(2\sqrt{\beta} \gamma),$$

with $\text{Im}(\beta) \leq 0$.

8To make manifest that this integral representation is real one can take advantage of the relation for $x \in \mathbb{R}$: $K_{\nu}(ix) = \frac{i}{\pi} K_{\nu}(ix)$.
then show how it can be evaluated to determine specific hit functions in different dimensions and of different orders.

C. An analytical example

We will now present an example where the analytical properties of the propagator can be used to show the convergence of the diagonal subsequence of Padé approximants. To this end we shall use the $\delta$-function potential $V(x) = \lambda\delta(x)$ in one dimension and consider propagation under its influence. It will be convenient to calculate the Fourier transform of the free propagator in one dimension,

$$\hat{K}_0(x, y; \omega) = \int_{-\infty}^{\infty} K_0(x, y; \tau)e^{i\omega\tau}d\tau = \frac{e^{-|x-y|\sqrt{\omega}}}{2\sqrt{\omega}} e^{-i\omega T}, \quad (30)$$

because the propagator difference, $K_1(x, y; \tau) - K_0(x, y; \tau) := \Delta K(\lambda)$, can be written in closed integral form based on $\hat{K}_0(x, y; \omega)$, a known result found (in Minkowski time), e.g., in [45]:

$$\Delta K(\lambda) = \frac{\lambda}{2\pi} \int_{-\infty}^{\infty} \hat{K}_0(x, 0; \omega)K_0(0, y; \omega)e^{-i\omega T}d\omega. \quad (31)$$

The above integral can be solved using contour integration—in particular the Heaviside step function $\theta(T)$ results from different contour choices for both signs of $T$. Doing so we arrive at the closed form

$$\Delta K(\lambda) = \frac{\lambda}{4} e^{\frac{\pi i}{4}} e^{\frac{\pi T}{\lambda}} \text{erfc} \left( \frac{\Delta + \lambda T}{2\sqrt{T}} \right) \theta(T). \quad (32)$$

Notice that $\Delta K(\lambda)$ is an entire function of $\lambda$. Furthermore, using the results we derive below for the hit functions in $D = 1$ dimension, the series (19) coincides exactly with the series expansion of $K = K_0 + \Delta K(\lambda)$ in powers of $\lambda$, which also informs us that it correctly produces the diagonal Padé approximants in this case. The analyticity now allows us to invoke a theorem due to Nuttall [46] to prove that the diagonal subsequence of Padé approximants to $\Delta K(\lambda)$ converges to $\Delta K(\lambda)$ almost everywhere in $\mathbb{C}$, so that convergence in the limits $\lambda \to \pm\infty$ is assured. Furthermore, in the large $\lambda$ limit we obtain the leading order contribution

$$\Delta K(\lambda \to \infty) \sim -\frac{e^{-\frac{\pi T}{\lambda}}}{\sqrt{4\pi T}} + O\left(\frac{1}{\lambda}\right). \quad (33)$$

which reproduces the result for the kernel in a half-space when $x$ and $y$ have the same sign. This kernel (but in $D = 3$ dimensions) is approximated using the Padé approximants of the conjecture in (16) in Sec. VI, where we show that the sequence of approximants (built by integrating $n$-hit functions over the plane) gives a good estimate of the propagator.

V. IDENTITIES AMONG $N$-HIT FUNCTIONS

It turns out that there are various identities relating the generalized hit functions for different spatial dimensions or number of intermediate points, $n$. In this section, for clarity we will write the space dimension $D$ explicitly in the $n$-hit function by adding a super-index in the form $\mathcal{H}_{n}^{(D)}$.

A. Changing the order $n$ and dimension $D$

We begin by considering hit functions defined in different dimension. By writing out (13) explicitly,

$$\mathcal{H}_{0}^{(D)}(z_1, \ldots, z_n|y, x; T) = \frac{1}{(4\pi)^{(n+1)/2}} \int_{0}^{T} d\tau_0 \cdots \int_{0}^{T} d\tau_n \prod_{i=1}^{n+1} \frac{e^{-\Delta_i^2/4(\tau_i - \tau_{i-1})}}{(\tau_i - \tau_{i-1})^{D/2}} \theta(T), \quad (34)$$

we see that the dependence on the dimension appears explicitly in the denominator, the prefactor and, implicitly in the definition of each $\Delta_i = |z_i - z_{i-1}|$ as a distance in $D$-dimensional Euclidean space. As a direct consequence of differentiation under the integral the following functional relation holds:

$$\mathcal{H}_{0}^{(D+2)}(z_1, \ldots, z_n|y, x; T) = \prod_{i=1}^{n+1} \left( \frac{-1}{2\pi D_i} \frac{\partial}{\partial D_i} \right) \mathcal{H}_{0}^{(D)}(z_1, \ldots, z_n|y, x; T) \quad (35)$$

$$:= D_n^+ \mathcal{H}_{0}^{(D)}(z_1, \ldots, z_n|y, x; T), \quad (36)$$

where we denote the operation of raising the dimension by $D_n^+$ that acts on a hit function of order $n$ (see below for its lowering partner). Here it is important to stress that (35) must be interpreted as a functional relation and not an equality; in particular the meaning of the $\Delta_i$ is different on both sides as they are to be calculated in spaces of dimensions $D + 2$ and $D$, respectively. Nevertheless this identity allows one to calculate hit functions in higher dimensions of the same parity once a given $n$-hit function is explicitly known. To illustrate this, specific examples are given in Sec. V B. It is also possible to write the inverse functional relation that will lower the dimension from $D \to D - 2$ by inverting, with the appropriate boundary conditions at $\Delta_i \to \infty$, the operation defined in (35):

$$\mathcal{H}_{0}^{(D-2)}(z_1, \ldots, z_n|y, x; T) = \prod_{i=1}^{n+1} \left( \frac{2\pi}{\Delta_i D_i} \int_{\Delta_i}^{\infty} d\Delta_i \right) \mathcal{H}_{0}^{(D)}(z_1, \ldots, z_n|y, x; T) \quad (37)$$

$$:= D_n^- \mathcal{H}_{0}^{(D)}(z_1, \ldots, z_n|y, x; T), \quad (38)$$

which defines the dimension-lowering operator $D_n^-$ of order $n$ (of course, in the form we have written these operations it is clear that $D_n^\pm$ depend upon the intermediate points, but the meaning of the abstract operators acting on the hit functions should be clear).

Turning now to relations between hit functions of different order, from (13) one immediately sees how to incorporate an additional intermediate point, $z_{er}$, say, into the hit function which allows us to define an order-raising operator $N_{er}^{(D)}(z_{er})$, that will insert the additional point $z_{er}$ in the $D$-dimensional hit function.
function according to\(^9\)

\[
\mathcal{H}_0^{(D)}(z_1, \ldots, z_{n-1}|y, x; T) = \int_0^T d\tau \mathcal{H}_0^{(D)}(z_1, \ldots, z_{n-1}|x, \tau) K_0(y, z_n; T - \tau) \quad (39)
\]


\[
:= N_D^\pm(z_n) \mathcal{H}_0^{(D)}(z_1, \ldots, z_{n-1}|y, x; T). \quad (40)
\]

This can also be expressed as a convolution product \([\mathcal{H}_0(z_1, \ldots, z_{n-1}|x, \tau) \ast K_0(y, z_n)|(T)\) which produces the required insertion point using the appropriate kernel. We can invert (39) by introducing a Dirac \(\delta\)-distribution under the integral with an operator, \(N_D^\pm(z_n)\), that will lower the order of the hit function by removing the point \(z_n\) by integrating it out (in \(D\) dimensions)

\[
\mathcal{H}_0^{(D)}(z_1, \ldots, z_{n-1}|y, x; T) = \frac{1}{2} \int d^Dz_n \left( \frac{\partial}{\partial T} - \nabla^2 \right) \mathcal{H}_0^{(D)}(z_1, \ldots, z_n|y, x; T) \quad (41)
\]


\[
:= N_D^\pm(z_n) \mathcal{H}_0^{(D)}(z_1, \ldots, z_{n-1}|y, x; T), \quad (42)
\]

where \(\nabla_y^2\) is the Laplacian acting only on the \(y\) coordinates. In this case this follows from the Schrödinger-like equation that the hit function satisfies that is inherited from the Schrödinger equation, a generalization of the Green function equation satisfied by the kernel itself:

\[
\frac{1}{2} \left( \frac{\partial}{\partial T} - \nabla^2 \right) \mathcal{H}_0^{(D)}(z_1, \ldots, z_n|y, x; T) = \delta^D(z_n - y) \mathcal{H}_0^{(D)}(z_1, \ldots, z_{n-1}|y, x; T). \quad (43)
\]

Consider now the special case \(D = 1\): Due to its exceptional simplicity there is another manner to relate the \(n\) to \(n+1\)-hit functions. To this end observe that for \(D = 1\) (29) becomes

\[
\mathcal{H}_0^{(1)}(z_1, \ldots, z_n|y, x; T) = \left( i \frac{\pi}{\Delta} \right)^n \int_{-\infty}^{\infty} d\nu e^{-\nu^2 A} \theta(T). \quad (44)
\]

Notice that the only dependence on \(n\) is in the prefactor and the power of \(p\) in the integrand whilst the dimensional dependence only enters implicitly through \(\Delta\). By differentiation we obtain

\[
\mathcal{H}_0^{(1)}(z_1, \ldots, z_{n-1}|y, x; T) = -2 \frac{\partial}{\partial \Delta} \mathcal{H}_0^{(1)}(z_1, \ldots, z_n|y, x; T), \quad (45)
\]

which must be interpreted, as above, as a functional relation considering that the meaning of \(\Delta\) is different on both sides of (45) on account of the change in the number of intermediate points; the inverse relation follows from considerations similar to those made for (37)

\[
\mathcal{H}_0^{(1)}(z_1, \ldots, z_n|y, x; T) = \frac{1}{2} \int_{-\infty}^{\infty} d\Delta \mathcal{H}_0^{(1)}(z_1, \ldots, z_{n-1}|y, x; T). \quad (46)
\]

\(^9\)The same identity holds if a potential \(U\) is present by replacing \(\mathcal{H}_0 \rightarrow H\) and \(K_0 \rightarrow K\); cf. the footnote in Sec. II.

The commutative diagram in Fig. 2 shows the mapping between the different \(n\)-hit functions.

As a consequence of the relations described we draw the conclusion that all \(n\)-hit functions are related in dimensions with the same parity, and in a given dimension different orders of hit function (values of \(n\)) are also related, as shown in the diagram in Fig. 2.

The relations shown in the commutative diagram, and described by (35), (37), (39), and (41) are general. In the following section we shall apply them to obtain closed formulas for the odd dimensional case on account of its special simplicity (since all \(n\)-hit functions can be related to the \(D = 1\) or \(D = 3\) case). In the even dimensional case we will provide simple integral representations for \(n > 1\). In the special case with \(D = 1\) spatial dimensions the additional identities (45) and (46) hold which help us to construct all hit functions as a basis for the odd dimensional case. We also provide some results in even dimensions (in particular \(D = 2\)).

**B. Examples**

Let us illustrate the use of the integral representation (29) and the relations derived in this section to determine the \(n\)-hit function for some cases of interest. In particular we shall focus on the results for \(D = 1\) and \(D = 3\) and find the hit function for arbitrary \(n\).

First, to recover the result (21), we use \(\nu = -\frac{1}{\sqrt{2}}\) and the well-known analytic expression for the Bessel functions, \(K_{\nu/2}(x) = \sqrt{\frac{\pi}{2x}} e^{-x/2} I_{\nu/2}(x)\). Thus in \(D = 3\) the integral representation, (29), of the \(n\)-hit function turns into

\[
\mathcal{H}_0^{(3)}(z_1, \ldots, z_n|y, x; T) = \frac{2}{(4\pi)^{3/4}} \prod_{i=1}^{n+1} \Delta_i \int_{-\infty}^{\infty} d\nu e^{-\nu^2 + i\nu \sum_{i=1}^{n+1} \Delta_i} \theta(T). \quad (47)
\]

Completing the square in the exponent and shifting integration variables results in a Gaussian integral that evaluates immediately to (21). We also note that this result is, of course, compatible with the convolution identity, (39). Indeed, using
(A3), we immediately verify

\[
N_3^+(z_n) H_0^{(3)}(z_1, \ldots, z_{n-1}|z_n; x, \tau) = \int_0^T d\tau \theta(T - \tau) \frac{e^{-i \omega T + i p \Lambda}}{[4\pi (T - \tau)]^{\frac{1}{2}}} \\
\times H_0^{(3)}(z_1, \ldots, z_{n-1}|z_n; x, \tau)
\]  
(48)

\[
= \frac{1}{4\pi} \frac{\Delta - \Delta_{n+1}}{\Delta_1 \Delta_2 \cdots \Delta_n} \int_0^T d\tau \frac{e^{-i \omega T + i p \Lambda}}{[4\pi (T - \tau)]^{\frac{1}{2}}} \\
\times \frac{e^{i (\alpha - \alpha_{n+1})^2}}{\tau^{3/2}} \theta(T)
\]  
(49)

\[
= H_0^{(3)}(z_1, \ldots, z_n|y, x; T).
\]  
(50)

Likewise, for the case \( D = 1 \) it is also possible, with a little more work, to arrive at a closed form result for the hit function of arbitrary order. We can use this to verify the formulas (35) and (37) for changing dimension. A key relation is (46) which defines a useful recursion relation between hit functions of different order. In this case, the special form of the one hit \((n = 1)\) function makes it possible to solve this relation in closed form.

Although we have already quoted the result for \( H_1^{(1)}(z|y, x; T) \) in (2) of the introduction, here we take advantage of the opportunity to recalculate it using the “master formula” derived in (29). Indeed, for \( D = 1 \) this reduces to (44) so with \( n = 1 \) we must calculate

\[
H_1^{(1)}(z|y, x; T) = \frac{i}{2} \int_{\partial \gamma - \infty}^{\partial \gamma + \infty} \frac{dp}{2\pi i} e^{-\pi T + ip\Lambda}.
\]  
(51)

It is clear that we must take care with the pole at \( p = 0 \) (this is already apparent in (44)). To deal with this we separate the integral as follows [this treatment is equivalent to using the prescription for the pole \( \frac{1}{p^{\frac{1}{2}}} = P(1/p) - i\pi \delta(p) \)]:

\[
H_0(z|y, x; T) = \frac{1}{2} \left[ -\int_{\gamma(\epsilon)} \frac{dp}{2\pi i} e^{-\pi T + ip\Lambda} \right]
\]  
(52)

where \( \gamma(\epsilon) \) is the semicircle of radius \( \epsilon \) above the pole at the origin; the second is a principal value integral that excludes the interval \((-\epsilon, \epsilon)\) and the limit \( \epsilon \to 0^+ \) is understood. The first integral is evaluated as half the contribution given by the residue theorem at the pole, while the second is a known integral that evaluates to an error function so we reproduce the prior result quoted in (2),

\[
H_0(z|y, x; T) = \frac{1}{4} \left[ 1 - \text{erf}(\frac{\Delta}{2\sqrt{T}}) \right] = \frac{1}{4} \text{erfc}(\frac{\Delta}{2\sqrt{T}}).
\]  
(53)

once we put \( m = \frac{1}{2} \) and recall the change in the normalization of the hit function adopted in this paper. It is now easy to verify (35) and (37) using (23) which involve two integrations and two differentiations respectively.

Using this result in the recursion relation (46) shows that the one-dimensional \( n \)-hit function can be expressed in terms of the well-known iterated integrals of the complementary error function (see, for example, [47, 48]) defined by

\[
\mathcal{J}^n\text{erfc}(z) := \int_0^\infty dz^2 \mathcal{J}^{n-1}\text{erfc}(z),
\]  
(54)

\[
\mathcal{J}^{-1}\text{erfc}(z) = \frac{2}{\sqrt{\pi}} e^{-z^2},
\]  
(55)

\[
\mathcal{J}^0\text{erfc}(z) = \text{erfc}(z).
\]  
(56)

With this notation we solve the recursion relation via

\[
H_0^{(1)}(z_1, \ldots, z_n|y, x; T) = \frac{1}{4} \left[ \int_0^\infty d\Delta \frac{\Delta^{n-1}}{\sqrt{T} \Delta} \text{erfc}\left(\frac{\Delta}{2\sqrt{T}}\right) \right],
\]  
(58)

\[
H_0^{(1)}(z_1, z_2|y, x; T) = \frac{T}{16} \left[ -2 \frac{\Delta}{2\sqrt{T} \sqrt{\pi}} e^{-\frac{\Delta^2}{4T} + 1} \text{erfc}\left(\frac{\Delta}{2\sqrt{T}}\right) \right],
\]  
(59)

in which, of course, \( \Delta \) is interpreted according to the order \( n \). Using the integral definition of \( \mathcal{J}^n\text{erfc} \) given above, and the general result for (21) in three dimensions it is a simple exercise in induction to prove that the identities (35) and (37) hold between dimension \( D = 1 \) and \( D = 3 \) for arbitrary order \( n \). For completeness we supply the functional form of these results for \( n = 2 \) and 3:

\[
H_0^{(1)}(z_1, z_2|y, x; T) = \frac{T}{16} \left[ -2 \frac{\Delta}{2\sqrt{T} \sqrt{\pi}} e^{-\frac{\Delta^2}{4T} + 1} \text{erfc}\left(\frac{\Delta}{2\sqrt{T}}\right) \right],
\]  
(59)

where \( \Delta \) is to be formed using two and three intermediate points respectively. Returning to the example given in Sec. IV C, in which \( V(x) = \lambda \delta(x) \), the integrals over the surface \( S \) in (16) select \( z_1 = 0 \) so that for all \( n \) we would have \( \Delta = |z_1 + |y| \). It is then easy to check that, thanks to the singular support of the potential, the series (16) recovers the expansion of the full kernel in powers of \( \lambda \) and as such the correct Padé approximants (as discussed above the analytic nature of the kernel implies that the diagonal approximants converge to the desired function).

Since the dimensional shift formula (35) allows us to use the one-dimensional hit functions to produce all hit functions in odd dimension, it is of interest to give more explicit formulas for the former. As illustrated in (58) and (59), we observe that integrating elements of the set \( \{ \text{erfc}(z), e^{-z^2} \} \) produces polynomials in \( z \) multiplied by these same elements. Thus this form of closure under integration implies that the hit function in one dimension is built from polynomials in \( \frac{\Delta}{2\sqrt{T}} \) multiplied by the error function or Gaussian exponential of the same variable. Combining the recursion
relation with the integral representation (44) and employing the generalization \( \frac{1}{(p+1)\pi} = \mathcal{P}(\frac{1}{p}) + i(\frac{1-p}{(p-1)!})\pi \delta^{(1-p)}(p) \) it is possible to determine the hit function in terms of special functions

\[
\mathcal{H}^{(1)}_{0}(z_{1}, \ldots, z_{n}|y, x; T) = \frac{e^{-\frac{\Delta^2}{4\sqrt{T}}}}{4\sqrt{\pi(n-1)!}} T^{\frac{n-1}{2}} \left[ \Gamma\left(\frac{n}{2}\right) F_{1}(\frac{n}{2} - 1, \Delta^2/4T) - (n-1) \frac{\Delta}{2\sqrt{T}} F_{1}(\frac{n-1}{2} - 1, \frac{\Delta^2}{4T}) \right],
\]

where \( F_{1}(a, b; z) \) is a confluent hypergeometric function of the first kind [44,47]. Of more practical use we can relate the hypergeometric functions and parabolic cylindrical functions and write the latter in terms of derivatives of the error function to arrive at the following helpful identities:

\[
F_{1}(\frac{n}{2}, 1 - \frac{1}{2}; z^2) = (-1)^{n-1} \frac{\Gamma\left(\frac{n+1}{2}\right)}{2(n-1)!} \frac{d^{n-1}}{dz^{n-1}} \left[ e^{z^2} \text{erfc}(z) + (-1)^{n-1} \text{erfc}(z) \right],
\]

\[
F_{1}(\frac{n+1}{2}, 3 - \frac{1}{2}; z^2) = (-1)^{n-1} \frac{\Gamma\left(\frac{4}{2}\right)}{4z(n-1)!} \frac{d^{n-1}}{dz^{n-1}} \left[ - e^{z^2} \text{erfc}(z) - (-1)^{n-1} \text{erfc}(z) \right].
\]

If we use these in (60), then it simplifies substantially to a derivative representation

\[
\mathcal{H}^{(1)}_{0}(z_{1}, \ldots, z_{n}|y, x; T) = \frac{1}{4(n-1)!} \frac{T}{4} \left[ \text{erfc}\left(\frac{\Delta}{2\sqrt{T}}\right) H_{n-1}\left(-i\Delta\right) \right]
\]

\[
\left. + \frac{2}{\sqrt{\pi}} e^{-\frac{\Delta^2}{4\sqrt{T}}} \sum_{k=1}^{n-1} \binom{n-1}{k} (-1)^{k} H_{n-k-1}\left(\frac{\Delta}{2\sqrt{T}}\right) H_{n-k-1}\left(-i\Delta\right) \right],
\]

which define the polynomial coefficients of the complementary error function and Gaussian function mentioned above. If we now take the second derivative of the above representation the recurrence relation for Hermite polynomials ensures that it satisfies the differential equation (for brevity we indicate the order as an argument)

\[
\Delta^2 \mathcal{H}^{(1)}_{0}(n) + 8T \Delta \mathcal{H}^{(1)}_{0}(n) - 8T n \mathcal{H}^{(1)}_{0}(n) = 0.
\]

The generating formula (63) is a Rodrigues formula for the solutions of (65). Finally, to show that this coincides with the recursion relation in (57) we note that it is easy to prove by induction the following recursion relation for the iterated integrals of the complementary error function:

\[
2n^3 \text{erfc}(z) + 2z^2 \text{erfc}(z) - 2z^2 \text{erfc}(z) = 0.
\]

which readily implies the differential equation (65).

With these explicit results it is now also a trivial exercise to confirm that the more general form of the operations that raise and lower the order, (39), (41), (45), and (46) are satisfied by the \( n \)-hit function in one dimension (which is important because we constructed these functions using alternative relations). For \( n \in \{1, 2, 3\} \) this can be verified on a case by case basis; the result for arbitrary order is again most easily proved by induction. It is also worth pointing out that the propagator for the \( \delta \)-function potential in Eq. (32) provides an alternative means of generating the hit functions for arguments \( z_{i} = 0 \):

The \( n \)-hit function corresponds to the coefficient of \( \lambda^{n} \) in its expansion about \( \lambda = 0 \). This concludes our determination of this is all we require.

For the case of even dimensions, the recurrence relations and dimensional shift formulas imply that it suffices to obtain the \( n = 1 \)-hit function in \( D = 2 \), say. To achieve this we invoke the integral representation (34) which for \( n = 1 \) and \( D = 2 \) becomes

\[
\mathcal{H}^{(2)}_{0}(z|y, x; T) = \frac{1}{(4\pi)^{2}} \int_{0}^{T} d\tau e^{-\frac{z_{1}^{2}}{\sqrt{T}} - \frac{z_{2}^{2}}{\sqrt{T}}} \theta(T)
\]

\[
= \frac{1}{8\pi^{2}} K_{0}\left(\frac{\Delta_{1}\Delta_{2}}{2T}\right) e^{\frac{z_{1}^{2}}{\sqrt{T}}} \theta(T),
\]

where we used (A4). Thus we may, in principle, produce the infinite towers of hit functions of arbitrary order in even dimension by pure calculus.

VI. NUMERICAL BENCHMARKING

To investigate the feasibility of the calculation of the propagator put forward by (16) we have chosen to test the proposal with specific examples in two different geometries, namely, the \( D = 3 \) spatial region \( G \) with boundary \( \partial G \), given either by a plane \( \mathbb{R}^{2} \) or by a sphere \( S^{2} \), where Dirichlet boundary conditions are imposed on the boundaries. For simplicity we fix the center of the sphere at the origin and the plane to be the \( x-y \) plane passing through the origin. In the case of the sphere the region \( G \) considered is the interior of the unit sphere.
while in the case of the plane the region $G$ consists of the half space in $\mathbb{R}^3$ bounded by the plane.

In order to carry out a meaningful test we must compare the predictions given by (16) against the exact analytical expressions for the propagators (found in Appendix C). This requires us to evaluate the $c_n$ up to a given order to determine a truncated approximation for the quotient of determinants. In practice we have calculated up to $n = 3$, which corresponds to truncating the perturbative expansion at order $\lambda^4$. Evaluating $c_n$ as given by (17) requires calculating multiple integrals over the boundary geometries that are in general difficult to determine analytically [except for some special cases; cf. (D6) and (D7)], especially for $n > 1$. As such we resorted to numerical (Riemann) integration over the boundaries in such cases, although we should point out that despite this the calculation of (16) remains semi-analytic in that numerical techniques only enter the calculation to evaluate the coefficients of (16) remains semi-analytic in that numerical techniques only enter the calculation to evaluate the coefficients $c_n$ defined in (17). In our calculations this evaluation was done for $n = 0, 1, 2, 3$, using the Gauss-Kronrod quadrature method already built into Mathematica 10.3.1 for the numerical integration over the surfaces $S$. We found that the relative error of this quadrature method is variable in our calculations but stays around $\epsilon_Q = 0.1\%$; we discuss the error estimation in Sec. D.3. The numerical evaluation of the integrals involved was greatly simplified by particular choices of coordinate system to avoid spurious singularities which are detailed in Appendix D.

To obtain a graphical representation of the comparison we have plotted the dependence on the transition time, $T$, of the difference $K(x, y; T) - K_0(x, y; T)$ while keeping the points $x, y$ fixed in different cases. The difference $K(x, y; T) - K_0(x, y; T)$, the so-called scattering or vacuum-subtracted propagator, is calculated in different manners for the purpose of comparison. First and foremost, we compare all our results against the prediction given by the analytical (exact) expressions, that we have plotted in black. Second, we plot the successive predictions, labeled $P_1^1, P_2^1, P_3^1$, for the vacuum-subtracted propagator given by formula (16) for $N = 1, 2, 3$ for the truncated perturbative series. To be explicit, the quotients in these cases evaluate to

\begin{align}
P_1^1(T) &= -\frac{c_0}{c_1}, \\
P_2^1(T) &= \frac{c_1^2 + c_0 c_3 - 2 c_0 c_1 c_2}{c_2^2 - c_1 c_3}, \\
P_3^1(T) &= -\frac{c_2^4 - 4 c_0 c_3 c_2 c_3 + c_1 c_3^2}{c_3^3}.
\end{align}

(68)

and will be plotted below as dashed green, orange, and brown lines, respectively. As will be seen the agreement we obtain is quite good already with just these three approximants, but one might object that improving the accuracy further would require many more coefficients $c_n$ to be calculated. While this is indeed a viable way to proceed we can in fact improve the estimates without determining more coefficients by exploiting a simple observation: It turns out that in both geometries the sequence of $P_N^1$ alternates about the exact value (black solid line), behavior that can be recognized even when the analytic result is unknown. We can use this to obtain a much better numerical estimate using convergence-acceleration methods.

As such we introduce the Shanks transformation as a tool to accelerate convergence in the $\{P_N^1\}$ sequence. The Shanks transformation is particularly well suited to hasten convergence in alternating sequences [49], and it is this behavior exhibited by $\{P_1^1, P_2^1, P_3^1\}$ that has prompted its application here. We call $S_1$ and $S_2$ the first and second iterated Shanks transformation of the sequence $\{P_N^1\}$ (see Appendix E), explicitly given by

$$S_1(T) = \frac{P_1^1 P_3^3 - (P_2^2)^2}{P_1^1 + P_3^3 - 2 P_2^2}, \quad S_2(T) = \frac{S_1 + P_1^1 - (P_2^2)^2}{S_1 + P_1^1 - 2P_2^2}.$$ 

(69)

In our plots we have displayed both $S_1$ and $S_2$ as solid blue and red lines, respectively, from which we can appreciate a noticeable improvement. As expected, $S_2$ (in red) gives the closest fit to the analytical value (in black). Another favorable consequence of the alternating behavior is that we can roughly estimate the systematic error using the Leibniz criterion discussed in Appendix D.3. We have included this estimation as error bars (in grey) centered around the $S_2$ curve. From a practical point of view we can consider $S_2$ to be the best estimate we obtained in both cases, using only the first four coefficients $c_n$ for its evaluation.

The plots we obtained are given in Figs. 3–5. They reveal that the accuracy achieved by this method is in most regions notably greater than our naïve estimate would suggest, even if this estimate remains a tighter bound where the deviation is largest.

A. Sphere

Our first case of study is the interior region of a unit sphere (a unit ball) bounded by the surface $S = \partial G \simeq S^2$. Two cases were investigated where the vacuum-subtracted propagator was calculated with the same initial point $x$ but distinct final points $y$. We called these two cases center-to-center propagation and center-off-center
propagation, corresponding to Cartesian coordinates \( x = y = (0, 0, 0) \) and \( x = (0, 0, 0), y = (0, 0, r) \), respectively.

We show the results obtained with numerical help for the multiple integration over the surface of the sphere in Fig. 3 and Fig. 4. In both plots we have avoided the point \( T = 0 \) (where the propagators are known to vanish), hence restricting to the interval \( 0.02 \leq T \leq 2.5 \). For larger values of \( T \) than those presented here we verified that the agreement improves asymptotically. In Fig. 4 we have taken the parameter \( r = 0.7 \), that is, \( y = (0, 0, 0.7) \) so that an appreciable deviation from the center-center propagation is observed.

For comparison to the exact result, the analytical value of the vacuum-subtracted propagator (in black) was calculated in both plots from the series representation developed in Appendix C [see (C7)] by truncating it at \( l_{\text{max}} = 3 \) and \( k_{\text{max}} = 8 \). These values for truncation are optimal because, in fact, adding more terms slows down computation while producing no noticeable effect in our plots. This can be seen by examining \( n_{4,1} \approx 8.18 \), so that the largest ignored term is numerically \( O(e^{-66.9T}) \) and the fast convergence rate of the series means the remaining terms do not contribute substantially. Other calculations are of course possible where \( x \) and \( y \) are completely arbitrary points (not necessarily colinear with the center of the sphere) and based on our numerical calculations we expect that a similar agreement would be found.

B. Plane

The same comparison was done for the region (half-space) bounded by an infinite plane, in this case \( S = \partial G \cong \mathbb{R}^2 \); we investigated propagation from a point \( p = (0, 0, d) \) to itself, called \( p \) to \( p \) propagation. An important difference with the previous section is that the plane is infinite, and hence in the numerical integration demanded by the calculation of the coefficients \( c_n \) a finite region \( G \) of the plane must be chosen that is large enough to introduce only small errors [this is possible due to the exponential fall off of the hit functions, (21), with distance from the endpoints]; this procedure is detailed in Appendix D.

Since the plane is infinite, in considering \( p \) to \( p \) propagation the only length scale is \( d \), the distance from \( p \) to the plane, therefore no generality is lost if we assume the numerical value \( d = 1 \) that we have chosen for simplicity in our numerical calculations. In Fig. 5 we chose for clarity the interval \( 0.1 \leq T \leq 6 \) since, as in previous plots, the agreement improves for larger values of \( T \). One can see that Fig. 5 shows how the results obtained for the plane have a comparable degree of accuracy to those found in the sphere. We interpret this good agreement as indicative that if further levels of accuracy are calculated in (16) one can obtain a sequence of approximants that will converge to the desired difference of propagators. Naturally, since we have considered only three different cases presented in Figs. 3–5 based on truncation of (16) we can only consider this as evidence to support our claimed conjecture.

VII. CONCLUSIONS

In this article we have generalized the so-called “hit function,” introduced in previous work, to a more general integral transform of the quantum mechanical propagator: This \( n \)-hit function gives the relative contribution to the kernel involving propagation between its arguments through \( n \) (the order) intermediate points and contains the full information of the of the quantum system in question, being indeed a kind of generalized Green function for the Schrödinger operator as illustrated by (43). The \( n \)-hit function is notionally defined as a constrained path integral but we have developed various formulas, including (Riemann) integral representations, of this hit function that prescribe how to determine it from knowledge of the kernel.

We have also presented two ways in which quantum mechanical propagators can be constructed from the general hit function. The first, Eq. (14), is the direct inverse transformation recovering the unconstrained kernel (that extends to the case that the hit function is computed for a nontrivial kernel, i.e., in the presence of a background potential). The second, outlined in (15) and developed in Secs. II and III, shows how the hit functions of different orders can be formed into a series representation of the propagator in localized potentials.
or Dirichlet boundary conditions on different geometries. For the case of Dirichlet boundary conditions, our method combines nicely with Padé extrapolation allowing us to access the strong coupling regime of the series representation with remarkable success. While this has been illustrated for the case of propagation in the presence of a conducting plate and in the interior of a conducting spherical shell where we were able to show close agreement with analytical results obtained by conventional methods, the approach presented here can be used to calculate numerically the propagator or heat kernel in a bounded region with an irregular shape where an analytical expression is not viable. The proposal we make here could be considered as outlining a form of perturbative expansion for application in the strong coupling regime—as illustrated in (16) this suggestion estimates the kernel as a quotient of matrix determinants involving matrices of increasing size, using the Padé method to extrapolate to large values of the coupling parameter so that in this sense this can be considered a nonperturbative method based on perturbation theory. We have also shown that the convergence to this region of parameter space can be sped up by taking advantage of the (iterated) Shanks transformation. In Sec. III B we point out that this proposal can easily be adapted to nonlocalized static potentials of more general shape. Although we have not provided a general proof of the scope and validity of the conjecture outlined by (16) and (19), as such an endeavor lies out of the purpose of the present work, we have provided an analytical example where this can be done in one dimension. 

As part of the development of this approach to calculating the kernel we have also presented some general properties of these hit functions and derived relationships between the functions in different dimensions and of different order. It is important to emphasize that these functional relations imply that knowledge of one hit function (for some order) in one even- and one odd-dimensional space is, in principle, sufficient to determine all hit functions. To this end we have supplied the explicit form of all n-hit functions in $D = 1$ (and in fact also for $D = 3$) and calculated the $1$-hit function in $D = 2$ from which all other hit functions can be determined. Of course, the difficult parts of the process may be in successfully applying the identities of Sec. V to obtain explicitly the hit functions of different orders in the required dimension and then compute the appropriate integrals for the geometry of the system in question.

Before finishing we consider future work and possible extensions of the results presented here. As discussed above, especially in the context of our applications to the planar and spherical geometries in Sec. VI, the generalized hit function finds a natural application to studies of the Casimir energies in different geometries based on the first quantized representation of the effective action, which make for difficult calculations using standard techniques except in simple or special cases. The results of this application will be published in a later communication for some Casimir geometries of phenomenological interest to estimate the vacuum energy density and Casimir force. We estimate that the numerical precision we have attained is enough for such applications at a relatively lower computational cost compared to Monte Carlo-style simulations [35].

Longer term, other extensions include the incorporation of internal angular momentum—where the hit function will distinguish the contribution to the kernel from propagating modes for different spin degrees of freedom—and the generalization to relativistic particles and the Minkowski space path integral which connects to quantum field theory via the worldline formalism (which would, of course, involve a Wick rotation of the 0 component of the trajectory embedding). The hit function approach can also be used in scattering problems where the effect of a localized potential on an incoming particle is to be accounted for, so we believe there is a possibility of using a relativistic generalization of the hit function to treat high-energy scattering problems in QFT in an alternative manner. We would also like to point out that scattering of electromagnetic waves can be investigated using the presented method. The main difficulty here will be to implement the correct boundary conditions on the vector potential $A(x, t)$ and the transversality of the electromagnetic fields. The Klein-Gordon propagator can already be obtained from the nonrelativistic Schrödinger propagator by means of the worldline formalism through an integral transform $[38,39]$. Since the approach presented here provides a pathway to nonperturbative calculations it would be of interest to find out the scope of validity of this technique.

Yet another application can be found if the hit function is interpreted within the framework of the heat or diffusion equation. Here the results of the present paper allow one to solve the heat equation or the diffusion equation in bounded regions numerically, a problem of present day relevance in the study of Brownian motion. We also hope that the method we have introduced might stimulate interest within the framework of heat-kernel techniques where it might be of some relevance in the calculation of Seeley-DeWitt coefficients $[41]$. 

In light of the success we have obtained in calculating point-to-point propagation amplitudes we might use this method to calculate the density of states $\rho(E) = \frac{1}{2} \text{Im} \int d^D x K(x, x; -iE)$ for a quantum particle trapped inside a cavity and investigate the effects of changes in the shape of the cavity. For a many particle quantum system, for instance a quantum gas at finite temperature $T$, a multiparticle generalization would require the use of the coordinate multivector $\{x\} = \{x_1, x_2, \ldots, x_N\}$; the method outlined here then would estimate the multiparticle point-to-point propagator $K(\{x\}, \{x\}; T)$ so that the partition function of the system could be calculated as $Z(\beta) = \int d^D x K(\{x\}, \{x\}; \beta)$. Its practical use will depend on whether the required integrations can be carried out.

Finally, we allow ourselves to make a more speculative observation regarding Feynman’s original construction of the path integral. The classic construction, as expounded in $[34]$, for example, uses the intuition supplied by the Young’s slits experiment to motivate dividing propagation into infinitesimal steps, integrating the intermediate points along the trajectory over all space. One could interpret the large $n$ limit of the generalized hit function we define here as a concrete realization of this splitting procedure, following which the continuum path integral could be defined as the $n \to \infty$ limit of (14). Of course, the precise nature in which this limit is determined requires further analysis to make this interpretation more
precise but as a first attempt one could consider reproducing the interference pattern for double slits of finite size using the \( n = 1 \) hit function integrated over a finite volume. It also follows from this analysis that the \( n \) hit function will be useful in investigating the quasiclassical evolution of a quantum system when the classical trajectory goes through the prescribed space-time points \((z_1, \tau_1), \ldots, (z_n, \tau_n)\).

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To prove (A2) we change variable to \( x = \frac{\sqrt{T}}{\sqrt{T - \tau}} \) which reduces the problem to (A2) to that of calculating (A1). Finally, by acting with \(-\frac{\sqrt{T}}{\sqrt{T - \tau}}\) we can relate (A2) to (A3). Particular cases of (A1), (A2), and (A3) are found in the Appendix of [34] when \( \alpha \) and \( \beta \) are purely imaginary. One may now solve the integrals in representation (13) by substituting (11) with \( D = 3 \) and repeated use of (A3). The key point is that \( D = 3 \) is special, since it is the only dimension where the \( \tau \) integral of the product of two free propagators is again proportional to a free propagator, this allows one to obtain directly

\[
\mathcal{H}_0(z_1, z_2, \ldots, z_n; y, x; T) = \frac{1}{(4\pi)^{\frac{3n}{2}}} e^{-\frac{\pi^2}{T^{3/2}}} \prod_{j=1}^{n} \left[ \Delta_j + 1 \left( \sum_{k=1}^{i} \Delta_k \right)^{-1} \right],
\]

which can be simplified down to (21).

**APPENDIX B: CONSTRAINED PATH INTEGRAL**

In this Appendix we provide two approaches to prove identity (10). One is based on a path decomposition induced by a natural proper time interval partition and the other based on the Fourier representation of Dirac \( \delta \)-distributions and generating function techniques.

**APPENDIX A: THREE-DIMENSIONAL N-HIT FUNCTION**

We outline the calculation of the closed form for the \( n \)hit function in \( D = 3 \). In what follows \( \alpha, \beta \in \mathbb{C} \) and \( \text{Re}(\alpha) \geq 0, \text{Re}(\beta) \geq 0 \). Consider first the following integral:

\[
\int_0^\infty \exp\left(-\frac{\sqrt{\alpha x^2 - \beta}}{x^2}\right) dx = \frac{1}{2} \sqrt{\frac{\pi}{\alpha}} e^{-\frac{\sqrt{\alpha}}{2}}. \quad (A1)
\]

This identity can be proved by changing to the variable \( y = \sqrt{\alpha x^2 - \beta} \) which renders (A1) in Gaussian form after using contour integration. Alternatively one may use a complex Cauchy-Schlömilch transformation in the integrand by changing variable to \( \xi = \frac{1}{2} \frac{\sqrt{x}}{\sqrt{\alpha}} \) and then use Cauchy’s theorem to deform the contour so that the \( \xi \) integration is shifted to run on the real axis; finally replacing \( \xi \to x \) and averaging both integrals casts the original integral into a Gaussian integral leading to the r.h.s. of (A1). This result extends the integral (3.325) in [44] to complex values of \( \alpha \) and \( \beta \). We have also found it useful to compute the following integrals:

\[
\int_0^T \exp\left(-\frac{\alpha}{T - \tau} - \beta \right) \frac{d\tau}{\sqrt{(T - \tau)^2}} = \frac{\sqrt{\pi}}{\alpha T} e^{-\frac{\sqrt{\alpha}}{T^3}}, \quad (A2)
\]

\[
\int_0^T \exp\left(-\frac{\alpha}{T - \tau} - \beta \right) \frac{d\tau}{\sqrt{\tau(T - \tau)^2}} = \frac{\sqrt{\pi}}{T^3} \frac{\sqrt{\alpha + \beta}}{\sqrt{\alpha \beta}} e^{-\frac{\sqrt{\alpha \beta}}{T^3}}, \quad (A3)
\]

\[
\int_0^T \exp\left(-\frac{\alpha}{T - \tau} - \beta \right) \frac{d\tau}{\tau(T - \tau)} = \frac{2}{T} K_0(2\sqrt{\alpha \beta}/T) e^{-(\alpha + \beta)/T}. \quad (A4)
\]

1. **First approach**

Given boundary conditions \( \chi(0) = x, \chi(T) = y \) a path integration with an action functional \( S(0, T) = \int_0^T L \, dt \) over all trajectories that satisfy these boundary conditions can always be split for any value \( s \) such that \( 0 < s < T \) as

\[
\int_{\chi(0) = x}^{\chi(T) = y} \mathcal{P} x e^{-S(0, T)} L \, dt = \int d^3 w \int_{\chi(0) = x}^{\chi(y) = y} \mathcal{P} x e^{-S(x, T)} L \, dt \int_{\chi(0) = x}^{\chi(T) = y} \mathcal{P} x e^{-S(x, T)}. \quad (B1)
\]

Consider now the identity found in [1,34] for \( 0 < \tau < T \)

\[
\int_{\chi(0) = x}^{\chi(T) = y} \mathcal{P} x e^{-S_0(0, T)} \delta^0(\chi(\tau) - z) = K_0(y, z; T - \tau) \times K_0(z, x; \tau), \quad (B2)
\]

where \( S_0(0, T) := \int_0^T dt' \frac{\mathcal{L}}{4} \). Given now a set of ordered proper times \( 0 \leq \tau_1 \leq \tau_2 \leq \cdots \leq \tau_n \leq T \) we can always find a set \( \mathcal{S} = \{ s_k : 0 \leq k \leq n \} \) so that \( 0 := s_0 \leq \tau_1 \leq s_1 \leq \tau_2 \leq s_2 \leq \cdots \leq s_{n-1} \leq \tau_n \leq s_n := T \) and evaluate \( \mathcal{I}_n \) by splitting the \( [0, T] \) interval precisely at each \( s_k \), that is,

\[
\mathcal{I}_n = \prod_{k=1}^n \int d^3 w_k \int_{\chi(0) = x}^{\chi(T) = y} \mathcal{P} x e^{-S_0(0, T)} \delta^0(\chi(\tau_k - s_k) - \chi), \quad (B3)
\]
where we defined \(w_0 := x, w_n := y\). Now use (B2) repeatedly to solve for each term to arrive at

\[
\mathcal{I}_n = \prod_{i=1}^n \int d^2 w_k K_0(w_k, z_k; \tau_k) K_0(z_k, w_{k-1}; \tau_k - s_{k-1}).
\]

(B4)

The integrations can now be performed directly using the reproducing kernel property; the result is independent of \(S\) and is in fact (10).

2. Second approach

Here we give a brief proof of the result (9) by explicit evaluation of the path integral. Using the Fourier decomposition of the \(\delta\)-functions under the integral we can write

\[
\mathcal{I}_n = \prod_{i=1}^n \int \left( \frac{d^2 k_i}{(2\pi)^2} e^{i\sum_{j=1}^i k_j z_j} \int_{(0)=x}^{x(T)=y} \mathcal{D}(\tau) e^{-\frac{\mathcal{D}}{T} \int [\dot{\mathcal{D}} + ij(\tau) \mathcal{D}(\tau)]} \right),
\]

(B5)

where \(j(\tau) = \sum_{i=1}^n k_i \delta(\tau - \tau_i)\). Completing the square in the path integral exponent and evaluation of the Gaussian integral lead to

\[
\mathcal{I}_n = (4\pi T)^{-n} \left( \frac{\mathcal{D}}{2T} \right)^n \prod_{i=1}^n \int d^2 k_i e^{i\sum_{j=1}^i k_j z_j} \int_{(0)=x}^{x(T)=y} \mathcal{D}(\tau) e^{-\frac{\mathcal{D}}{T} \int [\dot{\mathcal{D}} + ij(\tau) \mathcal{D}(\tau)]},
\]

(B6)

where \(\Delta_{ij} := \Delta(\tau_i, \tau_j) = \frac{1}{2} |\tau_i - \tau_j| - \frac{1}{2} (\tau_i + \tau_j) + \frac{\mathcal{D}}{2T}\) is a “worldline Green function” for the free particle subject to Dirichlet boundary conditions on \([0, T]\) and \(\hat{x}_i := \hat{x}(\tau_i) = x + (y - x) \frac{\tau_i}{T}\) is a parametrization of the straight line path between \(x\) and \(y\) in time \(T\).

The (Riemann) integrals over the \(k_i\) are also Gaussian and evaluate to give

\[
\mathcal{I}_n = (4\pi T)^{-n} e^{-\frac{\mathcal{D}^2}{2T}} (4\pi)^{-n/2} \left| \Delta_{1}^{1/2} \right| \prod_{i=1}^n \left( \frac{\mathcal{D}}{T} \right)^2 \left( \frac{\mathcal{D}}{T} \right)^{n-1} e^{-\frac{\mathcal{D}^2}{2T}} \sum_{\tau} e^{i\tau (\hat{x}_i - \hat{x}_j)}.
\]

(B7)

To confirm that this coincides with the product of propagators in (9) we fix the ascending order used throughout the main text. Then

\[
\text{det}(-\Delta) = \prod_{i=1}^n (\tau_i - \tau_{i-1}),
\]

(B8)

\[
\sum_{i,j=1}^n (z_i - \hat{x}_i) \Delta_{ij}^{-1} (z_j - \hat{x}_j) = \frac{(z_n - \hat{x}_n)^2}{T} - \sum_{i=1}^{n-1} \frac{(z_i - z_{i+1} - \hat{x}_{i+1} + \hat{x}_i)^2}{T} + \frac{(z_1 - \hat{x}_1)^2}{T},
\]

(B9)

Substituting these results into (B7) we recover the product in (10) for the chosen ordering. For evaluation of a similar path integral (essentially the Fourier transform of the one presented here) see [33,50,51].

APPENDIX C: EXACT PROPAGATORS

Here we show how the exact propagators were derived analytically for the interior region bounded by a spherical two-dimensional shell in a three-dimensional Euclidean space and for the half three-dimensional Euclidean space bounded by a two-plane; these propagators were obtained following the standard methods of mode summation or spectral decomposition and the image method, respectively.

1. Propagator inside the sphere

The propagator in question inside the unit sphere satisfies the Green equation

\[
\left( \frac{\partial}{\partial \tau} - \nabla^2 \right) K(x, x'; \tau) = \delta(\tau) \delta^3(x - x'),
\]

(C1)

along with the Dirichlet boundary conditions

\[
K(x, x'; 0) = \delta(x - x'), \quad K(x, x'; t) = 0, \quad \text{for } t > 0,
\]

\[
K(x, x'; t) = 0, \quad \text{for } |x| = 1.
\]

(C2)

We solve (C1) using separation of variables, which naturally leads to the following structure for the solution:

\[
K(x, x'; t) = \sum_{\mu, \ell, m} A_{\ell, m, \mu}(x') e^{-\mu t} J_\ell(\sqrt{T} r) Y_\ell Y_m \theta(t),
\]

(C3)

where \(\Omega = (\theta, \phi)\) are the spherical angle coordinates of \(x, r = |x|, J_\ell\) are the spherical Bessel functions, \(Y_{\ell m}\) are the spherical harmonics, and the coefficients \(A_{\ell, m, \mu}(x')\) are fixed by the boundary conditions (C2). First, we notice that necessarily \(\mu = u_{\ell, k}^2\) being \(u_{\ell, k}\) the zeros of \(J_\ell(r)\) for \(k = 1, 2, \ldots\), therefore the values of \(\mu\) are labeled by the index pair \((\ell, k)\). A way to proceed now is to observe that the orthogonality relation

\[
\int_0^1 x^2 j_k(x u_{\ell, k}) j_k(x u_{\ell, k}) \, dx = \frac{\delta_{kk'}}{2} [j_{k+1}(u_{\ell, k})]^2
\]

(C4)

implies in turn the completeness relation

\[
\frac{\delta(r - r')}{r^2} = 2 \sum_{k=1}^\infty \frac{j_k(r u_{\ell, k}) j_k(r' u_{\ell, k})}{j_{k+1}(u_{\ell, k})},
\]

(C5)
which taken together with the completeness relation for the spherical harmonics (where \( \cos \gamma = \hat{x} \cdot \hat{x}' \)),

\[
\sum_{\ell=0}^{\infty} \sum_{m=-\ell}^{\ell} Y_{\ell m}(\Omega')Y_{\ell m}(\Omega) = \frac{2 \ell + 1}{4\pi} P_\ell(\cos \gamma) = \delta(\Omega - \Omega'),
\]

and setting \( t = 0 \) in (C3) yield \( A_{\ell m, \mu}(x') = Y_{\ell m}(\Omega') j_\ell(r u_{\ell, \mu}) / j_{\ell+1}(u_{\ell, \mu}) \), finally

\[
K(x, x'; t) = \sum_{\ell=0}^{\infty} \sum_{m=-\ell}^{\ell} 2 \ell + 1 \frac{j_\ell(r u_{\ell, \mu})}{2\pi} \frac{j_{\ell+1}(u_{\ell, \mu})}{j_{\ell+1}(u_{\ell, \mu})} \times P_\ell(\cos \gamma) e^{-u_{\ell, \mu} t} \theta(t), \tag{C7}
\]

which can be cut off at appropriate values of \( l \) and \( k \).

2. Propagator in half space

For the simple geometry of Dirichlet boundary conditions along an infinite plane the analytic calculation of the kernel is trivialized by the method of images. Indeed, denoting by \( p' \) the point corresponding to the reflection of \( p \) in the plane, we may write the kernel on the half-space \( y \) the same side of the plane, as the image method requires, for \( y \) and \( x \) on different sides the Dirichlet boundary conditions imply \( K(y, x; T) = 0 \]

\[
K(y, x; T) = K_0(y, x; T) - K_0(y, x; T), \tag{C8}
\]

which, by construction, vanishes for points \( y = y^* \) (or, by symmetry of the kernel \( x = x' \)) on the surface of the plane and satisfies the appropriate Green equation on the half-space. In a path integral calculation, this result follows from noting that the contribution from trajectories defined in \( \mathbb{R}^2 \) that enter the forbidden region as they propagate to \( y \) are canceled by those that travel to the image point \( y^* \), leaving only those that remain in the correct half-space (including various reflections off the plane); see [52] for more information. Our calculation with the hit function will implement this cancellation by explicitly excluding paths that penetrate the boundary.

**APPENDIX D: SPHERE AND PLANE NUMERICS**

The main difficulty in the numerical evaluation of the coefficients \( c\alpha \) in (17) for the three-dimensional case is the presence of the denominator \( 1 / \prod_i \Delta_i \) in the integrand (17), since as this is integrated numerically it produces spurious divergences when adjacent points \( z_i \) and \( z_{i+1} \) become coincident. These divergences are only artificial and they are easily removed by using an appropriate coordinate system. This is in general the case; on a sufficiently regular surface \( S \) a coordinate system can always be defined so that all such spurious divergences disappear. We show how this is attained in our cases of study and give more details on the numerics that led to the results reported in Sec. VI.

1. Sphere

To eliminate the spurious divergences in \( c\alpha \) we introduce a coordinate system to evaluate (17), namely,

\[
(-1)^n \int_{\partial G} d\Omega_1 \cdots \int_{\partial G} d\Omega_n \mathcal{H}_0(z_1, \ldots, z_n; y, x; T), \tag{D1}
\]

where \( d\Omega_i \) is the solid angle measure for the sphere \( \partial G \) parametrized by coordinates for \( z_i \). These coordinates are specified as follows. First let the last point \( z_n \) have the standard \((\theta_n, \phi_n)\) spherical coordinates with origin at the center of the sphere, choosing the orientation for the \( z \) (azimuthal) axis so that the point \( y = (0, 0, r) \) lies along it. Next, consider the point \( z_i \) for \( 1 \leq i \leq n - 1 \) and assign it the spherical coordinates \((\theta_i, \phi_i)\). Again the origin of this \( i \)th coordinate system coincides with the center of the sphere but the azimuthal axis now goes through \( z_{i+1} \) so that the polar angle \( \theta_i \) is actually measured with respect to the radius that passes through \( z_i \).

This has been illustrated in Fig. 6 for \( n = 3 \), where the first choice is made for \( z_3 \) whose azimuthal axis is \( \hat{k}_3 \) (in red). Next, the axis is rotated to the direction \( \hat{k}_2 \) (in blue) and finally to \( \hat{k}_1 \) (in green). The polar coordinates of the point \( z_3 \) are described with respect to the polar system whose azimuthal axis is \( \hat{k}_i \). This peculiar choice of coordinates allows one to perform all the \( \phi_i \) integrations straightforwardly so that only the polar integrations over \( \theta_i \) remain to be done. These are greatly simplified through the change of variables

\[
\xi_i = \sin(\theta_i / 2), \quad 1 \leq i \leq n - 1, \quad \xi_n = \cos \theta_n. \tag{D2}
\]

In particular the lengths of segments and the total length of the polygonal path become

\[
\Delta_1 = 1, \quad \Delta_{n+1} = \sqrt{1 + r^2 - 2 r \xi_n}, \quad \Delta_i = 2 \xi_i, \quad 2 \leq i \leq n, \quad \Delta = 1 + \sqrt{1 + r^2 - 2 r \xi_n} + 2 \sum_{i=1}^{n-1} \xi_i. \tag{D3}
\]
In effect we can apply the following substitution rules when calculating (D1):

\[
\int \frac{d\Omega_i}{\Delta_{i+1}} \to 4\pi \int_0^1 d\xi_i, \quad 1 \leq i \leq n - 1,
\]

\[
\int d\Omega_n \to 2\pi \int_{-1}^1 d\xi_n.
\]  \hspace{1cm} (D4)

These rules simplify greatly the calculation, for example in the center-propagation \((r = 0)\) the \(n\)th coefficient becomes

\[
(-\lambda)^n \frac{2}{(4\pi T)^{3/2}} \int_0^1 d\xi_1 \cdots
\]

\[
\int_0^1 d\xi_{n-1} \left(1 + \sum_{i=1}^{n-1} \xi_i e^{-(1+\sum_{i=1}^{n-1} \xi_i^2)/T}\right),
\]  \hspace{1cm} (D5)

and the whole series can be calculated analytically as

\[
K(0, 0; T) = \frac{1}{(4\pi T)^{3/2}} - \frac{2}{(4\pi T)^{3/2}} e^{-1/T
\]

\[
+ \lambda^2 \frac{1}{(4\pi T)^{3/2}} T(e^{-1/T} - e^{-4/T}) + \cdots.
\]  \hspace{1cm} (D6)

Note that in the general case of center-off-center propagation the original \(2n\) integrations are reduced to \(n\).

2. Plane

In calculating the terms of the perturbative series (15) for the half space bounded by a plane the first two terms can be determined analytically in the case of propagation from a point \(p\) to itself after proper time \(T\) has elapsed

\[
K(p, p; T) = K_0(p, p; T) - \lambda \frac{1}{16\pi T} \text{erfc}(d/\sqrt{T}) + \lambda^2 \cdots,
\]  \hspace{1cm} (D7)

where \(d\) is the distance from \(p\) to the plane. There is no loss of generality in taking \(p = (0, 0, d)\) and the plane to be the \(x-y\) plane passing through the origin, numerically we take \(d = 1\).

Higher order terms require numerical integration to be evaluated, so to that end let us introduce polar coordinates when integrating over the plane with an origin coincident with the origin of the original coordinate system. As a first example consider the \(\lambda^2\) term in (D7), take \(z_1 = (\rho_1, \phi_1)\) so that the angle \(\phi_1\) is measured with respect to the \(x\) axis on the plane, then take \(z_2 = (\rho_2, \phi_2)\) where now \(\rho_2\) is the distance measured with respect to \(z_1\), and the angle \(\phi_2\) is measured with respect to radius \(\rho_1\) so that

\[
\Delta_1 = \sqrt{\rho_1^2 + \rho_2^2}, \quad \Delta_2 = \rho_2,
\]

\[
\Delta_3 = \sqrt{\rho_1^2 + \rho_2^2 - 2\rho_1\rho_2 \cos \phi_2 + d^2},
\]  \hspace{1cm} (D8)

and the measures on the plane become \(d\sigma_1 = \rho_1 d\rho_1 d\phi_1,\ d\sigma_2 = \rho_2 d\rho_2 d\phi_2\) this choice of coordinates removes the spurious singularity in the integrand \(1/\Delta\) as \(\rho_2 \to 0\) simiplifying the convergence of numerical integration. The coordinate \(\phi_1\) can be integrated over directly resulting in a \(2\pi\) factor and a well-behaved integrand. We found it convenient to introduce a cutoff in the radial distance being integrated so that in practice \(\rho_1, \rho_2 \in [0, \rho_{\text{max}}]\), with \(\rho_{\text{max}} = 20\) in calculating \(c_1\) and \(\rho_{\text{max}} = 10\) in calculating \(c_2\) and \(c_3\). Increasing the integration cutoff \(\rho_{\text{max}}\) produces numerically identical values within the quadrature accuracy:

\[
c_1 = \frac{(2\pi)}{(4\pi)^{3/2} T^{3/2}} \int_0^{\rho_{\text{max}}} d\rho_1 \int_0^{2\pi} d\phi_1 \int_0^{\rho_{\text{max}}} d\rho_2
\]

\[
\times e^{-d^2/4T} \rho_1 \Delta \rho_1 \Delta
\]

\[
\sqrt{(d^2 + \rho_1^2)(d^2 + \rho_2^2 + \rho_3^2 - 2\rho_1\rho_2\cos \phi_2)}.
\]  \hspace{1cm} (D9)

For \(n > 2\) we must introduce additional points \(z_k\). Now, the procedure is changed by affixing the origin of the polar coordinate system for point \(z_k\) at the point \(z_{k-1}\) and measuring the polar angle \(\phi_k\) always with respect to the \(x\) axis. For example, for \(n = 3\) the relevant cancellation occurs because \(\Delta_2 = \rho_2\) and \(\Delta_3 = \rho_3\), and again \(\phi_1\) can be integrated directly, resulting in (with \(\Delta = \Delta_1 + \Delta_2 + \Delta_3 + \Delta_4\))

\[
c_2 = -\frac{(2\pi)}{(4\pi)^{3/2} T^{3/2}} \int_0^{\rho_{\text{max}}} d\rho_1 \int_0^{\rho_{\text{max}}} d\rho_2 \int_0^{\rho_{\text{max}}} d\rho_3 \int_0^{2\pi} d\phi_1 \int_0^{2\pi} d\phi_2
\]

\[
\times e^{-d^2/4T} \rho_1 \Delta \rho_1 \rho_3 \Delta
\]

\[
\sqrt{\rho_1^2 + d^2}\Delta_4.
\]  \hspace{1cm} (D10)

where now

\[
\Delta_4 = \sqrt{\rho_1^2 + \rho_2^2 + \rho_3^2 - 2\rho_1\rho_2 \cos \phi_2 + 2\rho_1\rho_3 \cos \phi_3 + 2\rho_2\rho_3 \cos(\phi_2 - \phi_3)}.
\]  \hspace{1cm} (D11)

3. Error estimation

To assess the total error we made a simple estimate using the fact that the sequence of Padé-Shanks approximations for all our cases of study in seen to be alternating around the analytical answer (this behavior of course can be observed even if the analytical answer is unknown), a rough estimate is given by the Leibniz criterion as the difference of the last two approximations

\[
\epsilon(T) = |S_2(T) - S_1(T)|/|S_2(T)|.\]  \hspace{1cm} (D12)
APPENDIX E: PÉDÉ APPROXIMANTS AND SHANKS TRANSFORMATION

Pédé approximants are a sequence \[ \{P_N^M(x)\} \] of rational functions that approximate in a certain sense a given function \( f(x) \). In the main text this function is the quantum mechanical kernel. These rational functions are written as the quotient of polynomials \( P_N^M(x) = P_M(x)/Q_N(x) \) of orders \( M \) and \( N \), respectively. If the function \( f(x) \) is analytic and has a Taylor series \( f(x) = \sum_{n=0}^{\infty} c_n x^n \), then the functions \( P_N^M \) are constructed so that their Taylor series matches up to order \( M+N \) with that of \( f \):

\[
P_N^M(x) = c_0 + c_1 x + c_2 x^2 + \cdots + c_{M+N} x^{M+N} + O(x^{M+N+1}).
\]

(1.1)

This requirement, along with the normalization

\[
P_N^M(x) := \frac{\sum_{i=0}^{M} a_i x^i}{1 + \sum_{j=1}^{N} b_j x^j}
\]

(1.2)

uniquely determines the coefficients \( a_i \) and \( b_j \) through a linear system of equations that can be solved using Cramer’s rule; the result is well known in the literature [49]:

\[
P_N^M(x) = \begin{bmatrix}
x^N \phi_{M-N} & x^{N-1} \phi_{M-N+1} & \cdots & \phi_M \\
\phi_{M-N+1} & \phi_{M-N+2} & \cdots & \phi_{M+1} \\
\vdots & \vdots & \ddots & \vdots \\
\phi_M & \phi_{M+1} & \cdots & \phi_{M+N} \\
x^N & x^{N-1} & \cdots & 1 \\
\phi_{M-N+1} & \phi_{M-N+2} & \cdots & \phi_{M+1} \\
\vdots & \vdots & \ddots & \vdots \\
\phi_M & \phi_{M+1} & \cdots & \phi_{M+N}
\end{bmatrix}
\]

(1.3)

where \( \phi_L := \sum_{p=0}^{L} c_p x^p \). In the main text this method is applied to the Taylor expansion of the kernel given in (15).

In general if we are looking for the value \( f(\infty) \) the limit \( \lim_{n \to \infty} P_N^M(x) \) will be zero if \( N > M \) and divergent if \( N < M \), but the idea was put forward in [42] that \( M = N \) produces a finite limit that can be used to approximate \( f(\infty) \) for a certain class of functions \( f \). The ensuing subsequence of approximants \( P_N^N(x) \) are called the diagonal Pédé approximants. This is the underpinning idea in proposing (16), whose strong coupling limit we estimate by a sequence of Pédé approximants determined in terms of integrals of the hit function.

Shanks transformation

Given an infinite sequence \( \{a_n\} \) that converges to \( \alpha = \lim_{n \to \infty} a_n \) the convergence oftentimes behaves asymptotically in the large \( n \) limit as

\[ a_n \sim \alpha + Z q^n, \quad |q| < 1, \quad \text{as } n \to \infty, \]  

(1.4)

for some constants \( Z, q \). Given this hypothesis the first-order Shanks transformation [49] consists of the solution for a numerical estimator \( S \) for the limiting value of the sequence \( \alpha \), the idea is to replace \( \sim \) by \( = \) above and demand that the estimator \( S \) satisfies the linear system

\[ \alpha_{n+1} = S + Z q^{n+1}, \]  

(1.5)

\[ \alpha_n = S + Z q^n, \]  

(1.6)

\[ \alpha_{n-1} = S + Z q^{n-1}. \]  

The estimator found by solving this system is itself a sequence \( \{S(a_n)\} \) that relates nonlinearly to the original sequence \( \{a_n\} \) and is called the (first-order) Shanks transformation of \( \{a_n\} \):

\[ S(a_n) := \frac{\alpha_{n+1} a_{n-1} - \alpha_n^2}{\alpha_{n+1} - \alpha_n - 2\alpha_n}. \]

(1.6)

The Shanks transformation is useful because it often converges faster than the original sequence. When only a limited number of terms in the sequence are available it is often necessary to iterate the Shanks transformation to hasten convergence, for example, \( S(S(a_n)) \) is the second iteration. In the numerical calculations presented in Sec. VI we have used the first and second iteration, \( S_1 \) and \( S_2 \), to show how the numerical results obtained are already significantly improved. We point out that to use the Shanks transformation in this manner we need to calculate the coefficients at least up to the term \( \lambda^4 \) in (15). We should also point out that the validity of this method rests upon the hypothesis (E4) which can fail in some situations, it is nevertheless possible to lend support to the validity of this hypothesis if enough numerical data are available or if a theoretical arguments allow for its justification. Both the Pédé extrapolation method and the Shanks transformation-based acceleration of convergence are usually validated post hoc based on its rates of convergence and their numerical self-consistency; many more details can be found, e.g., in chapter 8 of [49].


