Lévy flights are random walks in which the probability distribution of the step sizes is fat-tailed. Lévy spatial diffusion has been observed for a collection of ultracold $^{87}$Rb atoms and single $^{24}$Mg$^+$ ions in an optical lattice, a system which allows for a unique degree of control of the dynamics. Using the semiclassical theory of Sisyphus cooling, we formulate the problem as a coupled Lévy walk, with strong correlations between the length $x$ and duration $\tau$ of the excursions. Interestingly, the problem is related to the area under the Bessel excursion. The latter describes the spatial diffusion from the origin without crossing it in between. Such a random walk in velocity space, which starts at the origin and does not cross it, describes the last jump event $x$ in the sequence. In the anomalous phases, the statistics of meanders and excursions are essential for the calculation of the mean-square displacement, indicating that our correction to the Montroll-Weiss equation is crucial and pointing to the sensitivity of the transport on a single jump event. Our work provides general relations between the statistics of velocity excursions and meanders on the one hand and the diffusivity on the other, both for normal and anomalous processes.

**I. INTRODUCTION**

The velocity $v(t)$ of a particle interacting with a heat bath exhibits stochastic behavior, which in many cases is difficult to evaluate. The position of the particle, assumed to start at the origin at time $t = 0$, is the time integral over the fluctuating velocity $x(t) = \int_0^t v(t')dt'$ and demands a probabilistic approach to determine its statistical properties. Luckily, the central-limit theorem makes it possible, for many processes, to predict a Gaussian shape for the distribution of this work. The position of the particle, assumed or model, is to compute $K_2$ and other transport coefficients. In normal cases, this can be done, at least in principle, via the Green-Kubo formalism, namely, by the calculation of the stationary velocity correlation function, which gives the diffusivity $K_2 = \int_0^\infty \langle v(t)v(0)\rangle dt$.

An alternative approach is investigated in this work and is based on the concept of excursions. We assume that the random process $v(t)$ is recurrent, and thus the velocity crosses the zero point, $v = 0$, many times in the observation window $(0, t)$. We divide the path $x(t)$ into a sum of increments,

$$x(t) = \int_0^{t_1} v(t') dt' + \int_{t_1}^{t_2} v(t') dt' + \cdots + \int_{t_{k-1}}^{t_k} v(t') dt' + \cdots$$

Here, $\{t_1, t_2, \ldots\}$ are the points in time of the velocity zero crossings, $v(t_i) = 0$. In the interval $(t_i, t_{i+1})$, the velocity is either strictly positive or negative. The velocity in each interval is thus a stochastic process that starts and ends on the origin without crossing it in between. Such a random curve is called an excursion. The random spatial increment

From the Area under the Bessel Excursion to Anomalous Diffusion of Cold Atoms

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\[ x_i = \int_{t_i}^{t_{i+1}} v(t') \mathrm{d}t' \] is the area under the excursion. The position of the particle, according to Eq. (1), is the sum of the random increments, namely, the sum of the signed areas under the velocity excursions, each having a random duration. The goal of this paper is to relate the statistics of the areas under these velocity excursions, and the corresponding random time intervals between zero crossings, to the problem of spatial diffusion. This connection is easy to find in the case in which the increments \( x_i \) and the duration of the excursions \( t_i = t_{i+1} - t_i \) are mutually uncorrelated, independent, and identically distributed random variables. Over a long measurement time, the number of excursions is \( t/\langle \tau \rangle \), where \( \langle \tau \rangle \) is the average time for an excursion. Then, according to Eq. (1), the mean-square displacement is \( \langle x^2 \rangle = \langle \chi^2 \rangle t/\langle \tau \rangle \), and hence we have

\[ K_2 = \frac{\langle \chi^2 \rangle}{2\langle \tau \rangle}. \tag{2} \]

This equation is reminiscent of the famous Einstein formula and shows that diffusion is related to the statistics of excursions. The original work of Einstein, discussed in many textbooks, is explicitly based on an underlying random-walk picture, and so it does not involve the area under random velocity excursions, nor zero crossings in velocity space. We will arrive at the simple equation (2) only at the end of our work, in Sec. XI, while here it merely serves as an appetizer to motivate the consideration of velocity excursions in some detail and to suggest the usefulness of developing tools for the calculation of \( \langle \chi^2 \rangle \) and \( \langle \tau \rangle \) (here, \( \langle \chi \rangle = 0 \), by symmetry). Obviously, Eq. (2) is based on the assumption that the variance \( \langle \chi^2 \rangle \) is finite and, as mentioned, that correlations are not important. The major effort of this paper is directed at considering a more challenging case, namely, a physically relevant stochastic process where the variance \( \langle \chi^2 \rangle \) diverges, and more importantly, the process exhibits correlations between \( \chi \) and \( \tau \). The particular system we investigate is a model for diffusion of atoms in optical lattices, following the experiments in Refs. [1–3] and the semiclassical theory of Sisyphus cooling [4–6].

One type of excursion that has been thoroughly investigated is the Brownian excursion [7–9]. A Brownian excursion is a conditioned one-dimensional Brownian motion \( v(t) \) over the time interval \( 0 < t < \tau \). The motion starts at \( v(0) = \epsilon \) (eventually, the \( \epsilon \to 0 \) limit is taken) and ends at \( v(\tau) = 0 \) and is constrained not to cross the origin \( v = 0 \) in the observation time \( 0, \tau \). The area under this curve is a random quantity whose statistical properties have been investigated by mathematicians [10–13]. More recently, this problem was treated with a path-integral approach and applied to the statistics of fluctuating interfaces [7,8,14,15], and related areas until the first passage time have been used to describe universality of sandpile models [16]. Here, we generalize the Brownian excursion to a process \( v(t) \) described by a Langevin equation, with an asymptotically logarithmic potential. We show how the random area under this Langevin excursion determines the dynamics of cold atoms. We believe that Langevin excursions, or more generally, random excursions, are useful tools in many areas of statistical physics; hence, their investigation beyond the well-studied Brownian excursion is worthwhile.

If \( v(t) \) is a Brownian motion or, as we will proceed to show, the velocity of cold atoms in shallow lattices, the set of points where \( v(t) = 0 \) has nontrivial properties, as it contains no isolated points and no interval. This is the case since we are dealing with a continuous path with power-law statistics of the crossing times (see details below). Mathematicians have investigated the statistics of level-crossings, e.g., zero-crossings, of continuous paths in great detail. The concept of local time, introduced by P. Lévy in the context of Brownian motion and Ito’s excursion theory for continuous paths, is a pillar in this field (see Refs. [17,18] and references within). Here, we use a heuristic approach, with a modification of the original process \( v(t) \), by introducing a cutoff \( \pm \epsilon \); the starting point of the particle after each zero hitting, which is taken to zero at the end. This expediency makes it possible to use renewal theory and continuous-time random walks (CTRW), which are tools well studied in the physics literature of discrete processes. In this sense, we differ from rigorous mathematical approaches. Thus, we avoid the problem of continuous paths, for example, the infinite number of zero crossings, by replacing the original path with an \( \epsilon \) modified path for which the number of zero crossings is finite (when \( \epsilon \) is finite). We show that physical quantities characterizing the entire process have a finite \( \epsilon \to 0 \) limit. For example, in Eq. (2), for Langevin dynamics of \( v(t) \), both \( \langle \chi \rangle \) and \( \langle \chi^2 \rangle \) approach zero as \( \epsilon \to 0 \); their ratio \( K_2 \) approaches a finite limit.

One might wonder why we wish to use excursions and their peculiar properties to evaluate the spreading of atoms or, more generally, other transport systems. The answer is that it turns into a useful strategy when the friction forces are nonlinear. In particular, we will investigate laser cooling, where within the semiclassical theory, the dimensionless friction force is \([5]\) (see details below)

\[ F(\nu) = -\frac{\nu}{1 + \nu^2}. \tag{3} \]

This friction force, induced by the laser fields, is linear for small velocities \( F(\nu) \sim -\nu \), similar to the Stokes friction law for a massive Brownian particle in water at room temperature. However, unlike such friction, which increases in magnitude with velocity, here, for large \( \nu \), \( F(\nu) \sim -1/\nu \to 0 \). Asymptotically, then, the system is frictionless. This implies that fast particles tend to remain fast for a long time, which in turn induces heavy-tailed
populations of fast particles [19]. In this case, as we show later, the standard picture of diffusion breaks down.

More specifically, the problem of diffusion of cold atoms under Sisyphus cooling was partially treated by Marksteiner et al. [6]. While clearly indicating the anomalous nature of the diffusion process, the main tool used was the evaluation of the stationary velocity correlation function of the process, followed by the use of the Green-Kubo formalism for the evaluation of $K_2$. They showed that for a certain critical value of the depth of the optical lattice, the value of $K_2$ diverges (see also Ref. [20]). Katori et al. [1] measured the mean-square displacement $\langle x^2 \rangle \sim t^{\theta_2}$ and recorded the onset of super-diffusion $\xi > 1/2$ below a critical depth of the optical lattice [1]. Wickenbrock et al. [2], in the context of driven optical lattice experiments, demonstrated with Monte Carlo simulations an upper limit on the spreading of the atoms $\langle x^2 \rangle \leq \text{const} \times t^\theta$. Sagi et al. [3] showed that a packet of spreading Rb atoms can be fitted with a Lévy distribution instead of the Gaussian distribution found for normal diffusion. These findings clearly indicate the breakdown of the usual strategy of treating normal diffusion and hence invite further theoretical investigations. The velocity correlation function is not stationary, and hence the Green-Kubo formalism must be replaced [21]. The moments exhibit multifractal behavior [22]; there is an enhanced sensitivity to the initial preparation of the system [23,24]: momentum fluctuations are described by an infinite covariant density [25]; and in certain parameter regimes, the Lévy central-limit theorem applies instead of the standard Gaussian version [6,26]. In this regime of shallow optical lattices, the power of the analysis of the area under Langevin excursions becomes essential, as we will show here. We note that the relation of laser cooling with Lévy statistics is not limited to the case of Sisyphus cooling considered in this manuscript. Subrecoil laser cooling, a setup different from ours, also leads to fundamental relations between statistical physics of rare events and laser-atom physics [27,28]. For a recent mini-review on the departure from Boltzmann-Gibbs statistical mechanics for cold atoms in optical lattices, see Ref. [29].

A. Scope and organization of paper

The current work significantly extends the investigation of the properties of the spatial distribution of atoms undergoing Sisyphus cooling that began in Ref. [26]. There we uncovered three phases of the motion which are controlled by the depth $U_0$ of the optical lattice: a Gaussian phase $x \sim t^{1/2}$, a Lévy phase $x \sim t^\theta$ and $1/2 < \theta < 3/2$, and a Richardson phase $x \sim t^{3/2}$ (see Eq. (80) below). Within the intermediate phase, the density of particles, in the central region of the packet, is described by a symmetric Lévy distribution, similar to the fitting function used to analyze a recent Weizmann Institute experiment [3]. However, this cannot be the whole story. As is well known, the variance of the Lévy distribution diverges, which implies that $\langle x^2 \rangle = \infty$, which is unphysical. Indeed, as mentioned, Katori et al. experimentally determined a finite mean-square displacement exhibiting super-diffusion (see also Ref. [2]). We showed related numerical evidence that the Lévy distribution is cut off at distances of the order $x \sim t^{1/2}$. This breakdown of Lévy statistics arises because of the importance of correlations between jump lengths $\chi$ and jump duration $\tau$, which are neglected in the derivation of the Lévy distribution. Our purpose here is to investigate these correlations in detail. As we proceed to show, the correlations are given by the statistical properties of the Langevin excursions discussed above, for all except the last interval. Because the velocity is not constrained to be zero at the measurement time $t$, this last interval is not described by an excursion but rather by a Langevin meander, where the random walk $v(t)$ begins at the origin and does not return for the entire duration of the walk. The properties of the excursions and meanders enter in a modified Montroll-Weiss [30] equation for the Fourier-Laplace transform of the density $P(x, t)$ that we derive. We then use this equation to calculate a quantity that is sensitive to the correlations, namely, the mean-square displacement. The mean-square displacement exhibits anomalous diffusion and is sensitive to the last jump event, i.e., to the statistics of the meander. Thus, our treatment modifies both the celebrated Montroll-Weiss equation, to include the last jump in the sequence (i.e., the meander), and the existing theory of areas under Brownian meanders and excursions, to include the dissipative friction force which is responsible for the cooling of the atoms. Our analysis illuminates the rich physical behavior and provides the needed set of mathematical tools beyond the decoupling approximation used to obtain the Lévy distribution in our previous work. For completeness, we present detailed decoupled and coupled analyses, the latter being the main focus of the current work.

The paper is organized as follows. We start with a brief survey of the semiclassical theory of Sisyphus cooling [5,6] and show the connection of the dynamics to Lévy walks following Marksteiner et al. [6]. The importance of the correlations between $\tau$ and $\chi$ is emphasized, a theme which, as mentioned, has not received its deserved attention. In Sec. IV, a simple scaling theory is presented, which yields the exponents describing the dynamics of the atomic packet. The main calculation of the distribution of the area under the Bessel excursion is found in Sec. V; the calculation of the area under the Bessel meander is given in Appendix E. A new coupled continuous-time random-walk theory in Sec. VI provides the connection between the statistics of excursions and meanders, and the evolution of the density profile. Asymptotic behaviors of the Fourier-Laplace transform of the joint probability density function (PDF) of jump lengths and waiting times is investigated in Sec. VII. These, in turn, give us the asymptotic behaviors of
the atomic density packet, the mean-square displacement, and the different phases of the dynamics that are investigated in Secs. VIII–XI. Derivation of the distribution of the time interval straddling time $t$ for the Bessel process is carried out in Appendix F, which allows us to connect between our heuristic renewal approach and more rigorous treatments [17,31] and further discuss the nontrivial fractal set of zero crossings.

II. SEMICLASSICAL DESCRIPTION OF COLD ATOMS—LANGEVIN DYNAMICS

We briefly present the semiclassical picture for the dynamics of the atoms. The trajectory of a single particle with mass $m$ is $x(t) = \int_0^t p(t')dt/m$, where $p(t)$ is its momentum. Within the standard picture [5,6,33] of Sisyphus cooling, two competing mechanisms describe the dynamics. The cooling force $F(p) = -\alpha p/[1 + (p/p_c)^2]$ acts to restore the momentum to the minimum energy state $p = 0$. Momentum diffusion is governed by a diffusion coefficient, which is momentum dependent, $D(p) = D_1 + D_2/[1 + (p/p_c)^2]$. The latter describes momentum fluctuations that lead to heating due to random emission events that stochastically jolt the atom. We use dimensionless units, time $t \rightarrow \tilde{t}a$, momentum $p \rightarrow p/p_c$, and distance $x \rightarrow x a/p_c$, and introduce the dimensionless momentum diffusion constant $D = D_1/(p_c)^2a$. For simplicity, we set $D_2 = 0$ since it does not modify the asymptotic $|p| \rightarrow \infty$ behavior of the diffusive heating term, nor that of the force, and therefore does not modify our main conclusions.

The Langevin equations

$$\frac{dp}{dt} = F(p) + \sqrt{2D\xi(t)}, \quad \frac{dx}{dt} = p$$

(4)

describe the dynamics in phase space. Here, the noise term is Gaussian, has zero mean, and is white, $\langle \xi(t)\xi(t') \rangle = \delta(t - t')$. The now dimensionless cooling force is

$$F(p) = -\frac{p}{1 + p^2}. \quad (5)$$

The force $F(p)$ is a very peculiar nonlinear friction force. We note that friction forces that decrease with increasing velocity or momentum like $1/p$ are found also for some nanoscale devices, e.g., an atomic tip on a surface [32].

One goal of the paper is to find the spatial distribution of particles governed by Eq. (4). The stochastic equation (4) gives the trajectories of the standard semiclassical model for the dynamics in an optical lattice, which in turn was derived from microscopic considerations [5,6]. Denoting the joint PDF of $(x,p)$ by $W(x,p,t)$, the Kramers equation reads

$$\frac{\partial W}{\partial t} + p \frac{\partial W}{\partial x} = \left[ D \frac{\partial^2}{\partial p^2} - \frac{\partial}{\partial p} F(p) \right] W. \quad (6)$$

From the semiclassical treatment of the interaction of the atoms with the counterpropagating laser beams, we have

$$D = cE_R/U_0, \quad (7)$$

where $U_0$ is the depth of the optical potential, $E_R$ is the recoil energy, and the dimensionless parameter $c$ [34] depends on the atomic transition involved [5,6,35]. $U_0$ is a control parameter; hence, different values of $D$ are attainable in experiment, and exploration of different phases of the dynamics are within reach [1–3,19]. Equation (7) is rather intuitive since deep optical lattices, i.e., large $U_0$, imply small $D$, while large recoil energy leads to a correspondingly large value of $D$.

The behavior of the distribution for momentum only (when $x$ is integrated out from the Kramers equation, yielding the Fokker-Planck equation) is much simpler and has been presented in previous work [25]. The equilibrium properties are governed by the effective potential in momentum space,

$$V(p) = -\int_0^p F(p')dp' = (1/2)\ln(1 + p^2), \quad (8)$$

which for large $p$ is logarithmic, $V(p) \sim \ln(p)$. This large $p$ behavior of $V(p)$ is responsible for several unusual equilibrium and nonequilibrium properties of the momentum distribution [1,19,23,25,36]. The equilibrium momentum distribution function is given by [19]

$$W_{eq}(p) = \frac{1}{Z}e^{-V(p)/D} = \frac{1}{Z}(1 + p^2)^{-1/2D}. \quad (9)$$

Here,

$$Z = \int_{-\infty}^{\infty} \exp\left(-\frac{V(p)}{D}\right)dp = \sqrt{\pi}T\left(\frac{1 - D}{2D}\right)\Gamma\left(\frac{1}{2D}\right). \quad (10)$$

is the normalizing partition function. Equation (9) is Student’s $t$ distribution, also sometimes called a Tsallis distribution [19]. Actually, the problem is, coincidentally, related to Tsallis statistics since, as mentioned, the equilibrium PDF is proportional to a Boltzmann-like factor, $W_{eq}(p) \propto \exp(-V(p)/D)$, so $D$ acts like a temperature. More importantly, the power-law tail of the equilibrium PDF, for sufficiently large $D$, implies a large population of fast particles, which in turn will spread in space faster than what one would expect using naive Gaussian central-limit theorem arguments. For example, if $1/D < 3$, the ensemble-averaged kinetic energy in equilibrium diverges since $\langle p^2 \rangle_{eq} = \infty$, while when $1/D < 1$, the partition function diverges and a steady-state equilibrium is never reached. A dramatic increase of the energy of atoms when the optical lattice parameter $U_0$ approaches a critical value was found experimentally in Ref. [1], and a power-law momentum distribution was measured in Ref. [19]. Of
course, the kinetic energy of a physical system cannot be infinite, and the momentum distribution must be treated as a time-dependent object within the infinite covariant density approach [25]. While much is known about the momentum distribution, both experimentally and theoretically, the experiments [1–3] demand a theory for the spatial spreading. We note that diffusion in logarithmic potentials has a vast number of applications (see Refs. [23–25,37–41] and references therein), e.g., Manning condensation [42], and unusual mathematical properties, including ergodicity breaking [36]. In general, logarithmic potentials play a special role in statistical physics [43–45].

III. MAPPING THE PROBLEM TO A LéVY WALK PROCESS

In principle, one may attempt to directly solve the Kramers equation (6) to find the joint PDF of the random variables \(x, p\) at a given time \(t\). In Fig. 1, we plot a histogram of the phase space obtained numerically. We see a complicated structure: Roughly along the diagonal, clear correlations between \(x\) and \(p\) are visible; on the other hand, along the \(x\) and \(p\) axes, decoupling between momentum and position is evident, together with broad (i.e., non-Gaussian) distributions of \(x\) and \(p\). At least to the naked eye, no simple scaling structure is found in \(x - p\) space, and hence, we shall turn to different, microscopic, variables that do exhibit simple scaling. This leads to an analysis centered on the mapping of the Langevin dynamics to a Lévy walk scheme [6,46] and the statistics of areas under random excursions.

Starting at the origin, \(p = 0\), the particle along its stochastic path in momentum space crosses \(p = 0\) many times (see Fig. 2). In other words, the random walk in momentum space is recurrent; this is the case even when \(D \rightarrow \infty\) since then the process in momentum space is one of pure diffusion (i.e., the force is negligible), and from Polya’s theorem, we know that such one-dimensional walks are recurrent. The cooling force, being attractive, clearly maintains this property.

Let \(\tau > 0\) be the random time between one crossing event and the next, and let \(-\infty < \chi < \infty\) be the random displacement for the corresponding \(\tau\). As shown schematically in Fig. 2, the process starting with zero momentum is defined by a sequence of jump durations \(\{\tau_1, \tau_2, \ldots\}\), with corresponding displacements \(\{\chi_1, \chi_2, \ldots\}\). These random waiting times and displacements generate a Lévy walk [6], as we explain presently. Here, \(\chi_1 = \int_{\tau_1}^{\tau_1 + \tau_2} p(\tau')d\tau'\), \(\chi_2 = \int_{\tau_1 + \tau_2}^{\tau_1 + \tau_2 + \tau_3} p(\tau')d\tau'\), etc. Let the points on the time axis \(\{\tau_1, \tau_2, \ldots\}\) denote times \(\tau_n > 0\) where the particle crosses the origin of momentum \(p = 0\) (see Fig. 2). These times are related to the waiting times: \(\tau_k = \sum_{i=1}^{k} \tau_i\). We see that \(\chi_k = \int_{\tau_k}^{\tau_k + \tau_{k+1}} p(\tau')d\tau'\) is the area under the random momentum curve constrained in such a way that \(p(\tau')\) in the time interval \((\tau_k, \tau_{k+1})\) does not cross the origin, while it started and ended there. The total displacement, namely, the position of the particle at time \(t\), is a sum of the individual displacements \(x = \sum_{i=1}^{n} \chi_i\). Here, \(n\) is the random number of crossings of zero momentum in the time interval \((0, t)\), and \(\chi^*\) is the displacement made in the last interval \((\tau_n, t)\). By definition, in the time interval \((\tau_n, t)\), no zero crossing was recorded. The time \(\tau^* = t - \tau_n\) is sometimes called the backward recurrence time (47). The measurement time is clearly \(t = \sum_{i=1}^{n} \tau_i + \tau^*\). In standard transport problems, the effect of the last displacement \(\chi^*\) on the position of particle \(x\) is negligible, and similarly, one usually assumes \(t = \tau_n\) when \(t\) is large. However, for anomalous diffusion of cold atoms, where the distributions

![Image](021036-5)
of displacements and jump durations are wide, these last events cannot be ignored.

One goal is to find the long-time behavior of \( P(x, t) \), the normalized PDF of the spatial position of a particle that started at the origin \( x = 0 \), \( p = 0 \) at \( t = 0 \). It is physically clear that, in the long-time limit, our results will not be changed if instead we consider narrow initial conditions, for example, Gaussian PDFs of initial position and momentum. Initial conditions with power-law tails will likely lead to very different behaviors \[23,24\]. Once we find \( P(x, t) \), we have the semiclassical approximation for the spatial density of particles. The latter can be compared with the Weizmann Sisyphus cooling experiments \[3\], provided that collisions among the atoms are negligible.

The Lévy walk process under investigation is a renewal process \[47\] in the sense that once the particle crosses the momentum origin, the process is renewed (since the Langevin dynamics is Markovian). This is crucial in our treatment, and it implies that the waiting times \( \tau_i \) are statistically independent, identically distributed random variables. Here, however, as we will soon discuss, the pairs \( \{\tau_i, x_i\} \) are correlated.

Since the underlying Langevin dynamics is continuous, we need a refined definition of the Lévy walk process. Both the \( \tau_i \)'s and the \( x_i \)'s are infinitesimal; however, the number of renewal events, \( n \), diverges for any finite measurement time \( t \), in such a way that the total displacement \( x \) is finite. In this sense, the Lévy walk process under investigation is different from previous works where the number of renewals, for finite measurement time, is finite. One way to treat the problem is to discretize the dynamics, as is necessary in any event to perform a computer simulation, and then \( x_i \) and \( \tau_i \) are, of course, finite. In our analytical treatment, following Marksteiner et al. \[6\], we consider the first passage time problem for a particle starting with momentum \( p_i \) and reaching \( p_f < p_i \) for the first time at \( \tau \). We take \( p_f = 0 \) and eventually \( p_i = e \to 0 \). The Lévy walk scheme is hence summarized with the following steps:

1. Choose, with probability \( 1/2 \), either \( +p_i \) or \( -p_i \).
2. Follow the Langevin dynamics until the particle reaches \( p_f = 0 \).
3. Record the random displacement \( x \) and random duration \( \tau \) during this excursion.
4. Go to step 1.

This loop is terminated at time \( t \), the final displacement \( x^* \) calculated, and as mentioned, the total displacement is \( x = \sum_{i=1}^{n} x_i + x^* \). In the first step, we have a probability of \( 1/2 \) to start with either \( +p_i \) or \( -p_i \) since the cooling force is antisymmetric and thus vanishes at \( p = 0 \). The advantage of presenting the problem as a set of recurrent random walks through the \( x \)'s and \( \tau \)'s instead of the direct Langevin picture stems from the fact that we can treat, analytically, the former Lévy walk picture.

We denote the joint PDF of the pair \( \{x, \tau\} \) of a single excursion by \( \psi(x, \tau) \). The theoretical development starts from the analysis of \( \psi(x, \tau) \), and then from this, we use the Lévy walk scheme to relate the single-excursion information to the properties of the entire walk and, in particular, \( P(x, t) \) for large \( t \).

There exists a strong correlation between the excursion duration \( \tau \) and the displacement length \( x \), which is encoded in \( \psi(x, \tau) \). Let the PDFs of \( x \) and \( \tau \) be denoted \( q(x) \) and \( g(\tau) \), respectively. Weak correlations would imply the decoupled scheme \( \psi(x, \tau) = q(x)g(\tau) \), which is far easier to analyze; however, as we shall see, this decoupling is generally not true, and it leads to wrong results for at least some observables of interest. Properties of \( q(x) \) and \( g(\tau) \) were investigated in Refs. \[6,35\] and are also studied below. The problem becomes more interesting and challenging because of these correlations. As we show, they are responsible for the finiteness of the moments of \( P(x, t) \), in particular, the mean-square displacement \( \langle x^2 \rangle \), and for the existence of a rapidly decaying tail of \( P(x, t) \). This, in turn, is related to the Lévy flight versus Lévy walk dilemma \[46\], to multifractality \[22\], and to the physical meaning of the fractional diffusion equation \[48\] used as a fitting tool for the Weizmann experiment \[3\] [see Eq. (105) and discussion therein]. As we show below, above the critical value \( D = 1 \), the correlations can never be ignored, and they govern the behavior of the entire packet, not only the large \( x \) tails of \( P(x, t) \). Physically, the correlations are obvious since long durations of flight \( \tau \) involve large momenta \( p \), which in turn induce large displacements \( x \). As an example of these correlations, we plot in Fig. 3 the displacement \( |x| \) versus the corresponding \( \tau \) obtained from computer simulation. The figure clearly demonstrates the correlations, and it also shows a \( |x| \propto \tau^{3/2} \) scaling, which we now investigate. Note how much simpler the \( x - \tau \) distribution is, compared to the \( x - p \) distribution in Fig. 1. Simulations presented in Fig. 3 were performed on a

![FIG. 3.](image-url) We plot the flight distance \( |x| \) versus the jump duration \( \tau \) to demonstrate the strong correlations between these two random variables. Here, \( D = 2/3 \), and the red line has a slope \( 3/2 \) reflecting the \( x \sim \tau^{3/2} \) scaling discussed in the text.
discrete lattice in momentum space, starting on the first lattice point (see Appendix D for details).

IV. SCALING THEORY—RELATION BETWEEN EXPONENTS

As shown by Marksteiner et al. [6] and Lutz [35], the PDFs of the excursion durations and displacements satisfy the asymptotic laws

\[ g(\tau) \sim g_1 \tau^{-3/2+\gamma}, \quad q(x) \sim q_1 x^{-(4/3+B)}, \] (11)

with

\[ \gamma = \frac{1}{2D}, \quad B = \frac{1}{3D}. \] (12)

In Appendixes A and B, we study these PDFs using backward Fokker-Planck equations. In particular, we find the amplitudes \( g_1 \) and \( q_1 \), and the relevant moments. The exponents \( \beta \) and \( \gamma \) are the most crucial aspects of these PDFs. When \( D \to \infty \), we get \( \gamma = 0 \), and Eq. (11) gives familiar limits. In the “high-temperature” limit of large \( D \), the cooling force is negligible, and the Langevin equation reduces to Brownian motion in momentum space. Then, \( g(\tau) \propto \tau^{-3/2} \), which is the well-known asymptotic behavior of the PDF of first passage times for unbounded one-dimensional Brownian motion [49,50]. Less well known is \( \lim_{D \to \infty} q(\tau) \propto |\tau|^{-4/3} \), which describes the distribution of the area under a Brownian motion until its first passage time (see Refs. [51,52], in which this PDF is given explicitly). Notice that the power-law behavior Eq. (11) yields a diverging second moment of the displacement \( q \) for \( D > 1/5 \), which in turn gives rise to anomalous statistics for \( x \).

The correlations between \( \chi \) and \( \tau \) are now related to the asymptotic behaviors of their PDFs, Eq. (11). We rewrite the joint PDF

\[ \psi(\chi, \tau) = g(\tau) p(\chi | \tau), \] (13)

where \( p(\chi | \tau) \) is the conditional PDF to find the jump length \( \chi \) for a given jump duration \( \tau \). We introduce a scaling ansatz, which is expected to be valid at large \( \tau \):

\[ p(\chi | \tau) \sim \tau^{-\eta} B \left( \frac{\chi}{\tau^\eta} \right). \] (14)

Since \( g(\tau) = \int_0^\infty \psi(\chi, \tau) d\chi \), we have the normalization condition \( \int_0^\infty B(\zeta) d\zeta = 1 \). By definition,

\[ q(\chi) = \int_0^\infty \psi(\chi, \tau) d\tau \propto \int_0^\infty d\tau \tau^{-3/2-\gamma-\eta} B \left( \frac{\chi}{\tau^\eta} \right). \] (15)

Changing variables to \( z = \chi/\tau^\eta \), we get

\[ q(\chi) \sim \text{const} \chi^{(1/3+\eta)/2}. \] (16)

Comparing with Eq. (11), we get a simple equation for the unknown exponent \( \eta \), which is \( 1 + 1/(2\eta) + \gamma/\eta = 4/3 + B \), so using Eq. (12), we find \( \eta = 3/2 \). This is precisely the scaling behavior \( \chi \sim \tau^{3/2} \) that we observe in our simulation, Fig. 3. Hence, the natural scaling solution of the problem is

\[ p(\chi | \tau) \sim \frac{1}{\sqrt{D} \tau^{3/2}} B \left( \frac{\chi}{\sqrt{D} \tau^{3/2}} \right). \] (17)

In hindsight, this result is related to Brownian scaling. For a particle free of the cooling force, its momentum will exhibit Brownian scaling, \( p \sim (D \tau)^{3/2} \), and hence the excursions that are integrals of the velocity scale as \( \chi \sim \sqrt{D} \tau^{3/2} \). The crucial point is that this simple Brownian scaling is maintained even in the presence of the cooling force for all \( D \). This is because of the marginal nature of the weak \( 1/p \) friction force for large \( p \), which leaves the Brownian scaling intact but changes the scaling function. While the \( 3/2 \) scaling is simple and \( D \) independent, the shape of \( B(\cdot) \) is sensitive to this control parameter. In the next section, we investigate \( B(\cdot) \), and for this, we must go beyond simple scaling arguments.

V. AREA UNDER THE BESSEL EXCURSION

A natural generalization of the Brownian excursion is a Langevin excursion. Such a stochastic curve is the path \( \chi(\tau) \), given by the Langevin equation, in the time interval \( 0 \leq t \leq \tau \), such that it starts and ends at \( p_i = p_f = \epsilon \) but is constrained to remain positive in between. Here, \( p_i = p(0) \) is the initial, and \( p_f = p(\tau) \) is the final location in momentum space. For our application, the path is considered in the limit \( \epsilon \to 0 \). Since the path never crosses the origin, the area under such a curve is \( A(\tau) = \int_0^\tau p(\tau') d\tau' \), and hence the PDF of \( A \) for fixed \( \tau \) yields the sought-after conditional PDF \( p(\chi | \tau) \). Obviously, \( \chi \) is the integral over the constrained path \( p(\tau') \); hence, by definition, it is the area under the excursion. The meander will describe the last jump \( \chi^* \) since at the measurement time the particle’s velocity is generally nonzero. For now, we will discuss only excursions, find \( p(\chi | \tau) \), and return to the meander later.

Here, we focus our attention on a specific excursion. We call the Bessel excursion, corresponding to the case \( F(\chi) = -1/p: \)

\[ \frac{dp}{dt} = -1/p + \sqrt{2D} \xi(t), \] (18)

so that the effective potential is the nonregularized logarithm, \( V(p) = \ln(p) \) and \( p > 0 \). Since the scaling approach is valid for long times, where excursions are long, the typical momentum \( p \) is large, and the details of the force
field close to the origin are negligible for the purpose of the calculation of the scaling function \(B(\cdot)\). We will check this assumption with numerical simulations, which of course use a regularized form of the force law. Some sample paths of Bessel excursion are presented in Fig. 4. The name Bessel excursion stems from the fact that Langevin dynamics in the nonregularized friction force field \(1/p\) corresponds to a well-known stochastic process called the Bessel process [41,53]. More information on the regularized and nonregularized processes is given in Appendixes A and B.

From symmetry, \(p(\chi|\tau) = p(-\chi|\tau)\), and so we may restrict our discussion to \(\tau > 0\) and \(\chi > 0\). Later, we will use this symmetry, which implies that for \(-\infty < \chi < \infty\) we have \(p(\chi|\tau) = p(|\chi|, \tau)/2\), where \(p(|\chi|, \tau)\) is the PDF of positive excursions, normalized according to \(\int_{0}^{\infty} p(|\chi|, \tau) d|\chi| = 1\).

Let \(G_{s}(x, p|p_{f})\) be the joint PDF of the random variables \(x\) and \(p\). Since \(\int_{0}^{\infty} p(\chi|\tau) d\chi\), we have \(\chi = 0\) at time \(\tau = 0\), so initially \(G_{s=0}(\chi, p|p_{f}) = \delta(p - p_{f})\delta(x)\). Later, we will take \(p\) to be the final momentum \(p_{f}\), which similarly to \(p_{i}\), will be set to the value \(\epsilon \to 0\) (for the sake of notational brevity, we omit the subscript in \(p_{f}\)). The calculation of \(p(\chi|\tau)\) follows three steps. For the Brownian case \(F(p) = 0\), this method was successfully applied by Majumdar and Comtet [8].

(6) We find the Laplace \(\chi \to sD\) transform of \(G_{s}(x, p|p_{f})\),

\[
\int_{0}^{\infty} e^{-sD\chi} G_{s}(x, p|p_{f}) d\chi = \hat{G}_{s}(s, p|p_{f}). \tag{19}
\]

The reason why we multiply \(s\) by \(D\) in the Laplace transform will become clear soon. Since \(\chi\) is a functional of the path \(p(\cdot)\), we will use the Feynman-Kac (FK) formalism to find \(\hat{G}_{s}(s, p|p_{f})\) (see details below). The constraint that the path \(p(\cdot)\) is always positive enters as an absorbing boundary condition at \(p = 0\) [41].

(ii) The second step is to consider \(p_{i} = p = \epsilon \to 0\) and obtain the Laplace transform

\[
\hat{p}(s|\tau) = \int_{0}^{\infty} p(\chi|\tau)e^{-sD\chi} d\chi = \lim_{\epsilon \to 0} \frac{\hat{G}_{s}(s, \epsilon|\epsilon)}{\epsilon}. \tag{20}
\]

The denominator in the above equation ensures the normalization since we must have \(\hat{p}(s|\tau)|_{\epsilon=0} = 1\).

(iii) Finally, we invert Eq. (20) back to \(\chi\) space using the inverse Laplace transform, \(s \to \chi\), which yields \(p(\chi|\tau)\). From there, we can also immediately find the scaling function \(B(\cdot)\), Eq. (14).

We now implement these steps to solve the problem.

The FK formalism [52] treats functionals of Brownian motion. Here, we use a modified version of the FK equation to treat overdamped Langevin paths [54]. Let \(A = \int_{0}^{\infty} U(p|\tau) d\tau\) be a functional of the Langevin path and assume \(U(\cdot) > 0\) so that we are treating positive functionals. Here, \(G_{s}(A, p|p_{f})\) is the joint PDF of \(A\) and \(p\), and \(\hat{G}_{s}(s, p|p_{f})\) is the corresponding Laplace \(A \to sD\) transform. The generalized FK equation reads

\[
\frac{\partial \hat{G}_{s}(s, p|p_{f})}{\partial \tau} = \left[\hat{L}_{D} - sD U(p)\right] \hat{G}_{s}(s, p|p_{f}). \tag{21}
\]

Here, \(\hat{L}_{D} = D(\partial p)^{2} - \partial p F(p)\) is the Fokker-Planck operator. When \(F(p) = 0\), Eq. (21) is the celebrated FK equation, which is an imaginary-time Schrödinger equation, and \(-sD U(p)\) is the potential of the corresponding quantum problem. The constraint on positive excursions, namely, \(p > 0\), gives the boundary condition \(\hat{G}_{s}(s, 0|p_{f}) = 0\). In the quantum language, this is an infinite
potential barrier for \( p < 0 \). This formalism can be used, in principle, to obtain the area under Langevin excursions for all forms of \( F(p) \).

For the Bessel excursion under investigation here, the functional \( U[p(t)] = p(t) \) is linear since \( x = \int_0^t p(t')dt' \). Quantum mechanically, this gives a linear potential and hence the connection to the Airy function found for \( F(p) = 0 \) in Refs. [7,8]. With the force field \( F(p) = -1/p \), we have \( \dot{L}_0 = D(\partial_p)^2 + \partial_p p^{-1} \), and hence we find, using Eqs. (18), (19), and (21) (the former gives the \( sD \) term),

\[
\frac{\partial G_i(s, p|p_i)}{\partial r} = \left( D \frac{\partial^2}{\partial p^2} + \frac{1}{D} \frac{\partial}{\partial p} - sDp \right) G_i(s, p|p_i).
\]

(22)

The solution of this equation is found using the separation ansatz

\[
G_i(s, p|p_i) = \sum_{k=0}^{\infty} a_k \psi_k(p)e^{-DE_i t},
\]

(23)

which yields the time-independent equation

\[
\frac{\partial^2}{\partial p^2} \psi_k + \frac{1}{D} \frac{\partial}{\partial p} \psi_k - sD\psi_k = -E_i \psi_k.
\]

(24)

We now switch to the more familiar one-dimensional Schrödinger equation via the similarity transformation \( \phi_k = p^{1/(2D)} \psi_k \), which gives

\[
-\frac{\partial^2}{\partial p^2} \phi_k + \left[ \frac{1}{2D} \left( \frac{1}{2D} + 1 \right) p^{-2} + sp \right] \phi_k = E_i \phi_k.
\]

(25)

This Schrödinger equation has a binding potential, which yields discrete eigenvalues, with an effective repulsive potential with a \( p^{-2} \) divergence for \( p \to 0 \) and a binding linear potential for large \( p \) provided \( s \neq 0 \). Equation (25) describes a three-dimensional nonrelativistic quantum particle in a linear potential [55,56], the \( p^{-2} \) part corresponding to an angular-momentum term.

As usual, \( \phi_k(p) \) yields a complete orthonormal basis \( \int_0^\infty \phi_k(p)\phi_m(p)dp = \delta_{km} \), and the formal solution of the problem is

\[
G_i(s, p|p_i) = \sum_k \left( \frac{p_i}{p} \right)^{1/(2D)} \phi_k(p_i)\phi_k(p)e^{-DE_i t}.
\]

(26)

We can scale out the Laplace variable \( s \) from the time-independent problem, defining \( \phi_k(p) = s^{1/6}f_k(s^{1/3}p) \) and \( E_k = s^{2/3}\lambda_k \). Using Eq. (25), we find

\[
-\frac{\partial^2}{\partial p^2} f_k + \left[ \frac{1}{2D} \left( \frac{1}{2D} + 1 \right) \tilde{p}^{-2} + \tilde{p} \right] f_k = \lambda_k f_k.
\]

(27)

where \( \tilde{p} = s^{1/3}p \). Notice that when \( D \to \infty \), the second term vanishes and we get the Airy equation. For this technical reason, we have chosen in Eq. (19) the Laplace variable as \( sD \) not \( s \). Also, it should be noted that the \( f_k \) are orthonormal in the variable \( \tilde{p} \).

It follows from Eq. (27) and the absorbing boundary condition that

\[
f_k(\tilde{p}) \sim \tilde{d}_k \tilde{p}^{1+1/(2D)}, \quad \tilde{p} \to 0,
\]

(28)

with \( \tilde{d}_k \) a \( k \)-dependent coefficient. The \( \tilde{d}_k \) are evaluated from solutions of Eq. (27) with the normalization condition \( \int_0^\infty (f_k(p))^2 dp = 1 \). For the initial and final conditions under investigation, \( p_i = p = \varepsilon \), we have

\[
\hat{G}_i(s, \varepsilon|e) \sim s^{2/3}\alpha_1^{1+2a} \sum_k (\tilde{d}_k)^2 e^{-D_k s^{2/3}},
\]

(29)

where the exponents \( \nu \) and \( \alpha \) will turn out to be useful,

\[
\nu = \frac{1 + D}{3D}, \quad \alpha = \frac{1 + D}{2D}.
\]

(30)

Note that classification of boundaries for a nonregularized Bessel process was carried out in Ref. [41] and is discussed here in Sec. XII. For an absorbing boundary condition, both the sign of the probability current and the usual condition of the vanishing of the probability on the absorbing point must be taken into consideration [41].

According to Eq. (20), we need \( \hat{G}_i(s = 0, e|e) \), in order to normalize the solution. This \( s = 0 \) propagator can be found exactly. The eigenvalue problem now reads

\[
-\frac{\partial^2}{\partial p^2} \phi_0 + \left( \frac{1}{2D} \left( \frac{1}{2D} + 1 \right) \right) \frac{\phi_0}{p^2} = E_0 \phi_0.
\]

(31)

The superscript 0 indicates the \( s = 0 \) case. Since \( s = 0 \), the linear field in Eq. (25) is now absent, so the “particle” is not bounded; hence, one finds a continuous spectrum \( E_0 = k^2 \). The wave functions consistent with the boundary condition are

\[
\phi_0^k(p) = B_k \sqrt{p} J_\alpha(kp),
\]

(32)

where \( J_\alpha(\cdot) \) is the Bessel function of the first kind. The second solution with \( J_{\alpha-1}(kp) \) is unphysical because of the boundary condition [41]. The normalization condition is

\[
1 = (B_k)^2 \int_0^L p(J_\alpha(kp))^2 dp,
\]

(33)

where \( L \to \infty \) is the “box” size that will eventually drop out of the calculation. Solving the integral in Eq. (33) gives \( B_k = \sqrt{\pi k/L} \). For initial and final conditions \( p_f = p_i = \varepsilon \), the \( s = 0 \) propagator then reads
\[ \mathcal{G}_s(z = 0, e|e) = \frac{re}{L} \sum_{k=0}^{\infty} k[J_{\alpha}(ke)]^2 e^{-Dk^2\tau}. \]  

(34)

Using the small \( z \) expansion of the Bessel function \( J_\alpha(z) \sim (z/2)^\alpha/\Gamma(1 + \alpha) \), one finds

\[ \mathcal{G}_s(z = 0, e|e) \sim \frac{e^{z+2\alpha}}{2^{2\alpha}\Gamma(1 + \alpha)} \sum_k k^{2\alpha+1} e^{-Dk^2\tau}. \]  

(35)

Notice that replacing \( J_\alpha(\cdot) \) with \( J_{-\alpha}(\cdot) \) leads to \( G_s(z = 0, e|e) \sim e^{z-2\alpha} \), which diverges since \( \alpha > 1/2 \); as mentioned, this unphysical solution is rejected because of the absorbing boundary conditions [41]. The usual density of states calculation follows the quantization rule of a particle in a box extending between 0 and \( L \), which gives \( kL = n\pi \) with integer \( n \). Hence, \( \sum_{k = -\infty}^{\infty} \frac{d}{\tau} \frac{d}{\tau} \), and one finds, after a change of variable of \( x = k^2 \),

\[ \mathcal{G}_s(z = 0, e|e) = \frac{e^{z+2\alpha}}{2^{2\alpha}\Gamma(1 + \alpha)} \int_0^{\infty} dx x^{\alpha-1} e^{-Dx}. \]

(36)

Inserting Eqs. (29) and (36) into Eq. (20), we find our first main result:

\[ \hat{p}(s|\tau) = \frac{2^{2\alpha+1}\Gamma(1 + \alpha)}{\xi(D\tau)^{3/2}} \sum_{n=0}^{\infty} \hat{d}_k^n e^{-D\lambda_k s^{2\alpha+1}/\tau}. \]  

(37)

In the limit \( D \rightarrow \infty \), we find \( |\hat{d}_k| \rightarrow 1 \), and the \( \lambda_k \)'s are the energy eigenvalues of the Airy equation, related to the zeros of the Airy function; up to a rescaling of \( s \), we get the result obtained by Darling [10] and Louchard [11] for the area under the Brownian excursion, namely, the Laplace transform of the Airy distribution [8]. We can show that \( \hat{p}(s = 0|\tau) = 1 \), as expected.

TABLE I. Table of the first five eigenvalues and coefficients \( \hat{d}_k \), for \( D = 2/5 \).

<table>
<thead>
<tr>
<th>( k )</th>
<th>( \hat{d}_k )</th>
<th>( \lambda_k )</th>
</tr>
</thead>
<tbody>
<tr>
<td>0</td>
<td>0.41584</td>
<td>3.5930</td>
</tr>
<tr>
<td>1</td>
<td>-0.56973</td>
<td>5.0753</td>
</tr>
<tr>
<td>2</td>
<td>0.68170</td>
<td>6.3754</td>
</tr>
<tr>
<td>3</td>
<td>-0.77287</td>
<td>7.5582</td>
</tr>
<tr>
<td>4</td>
<td>0.85117</td>
<td>8.6567</td>
</tr>
</tbody>
</table>

Each term on the right-hand side can be inverse transformed separately since the inverse Laplace transform of \( (Ds)\gamma \) is \( x^{\gamma-1}/\Gamma(-\gamma) \). After term-by-term transformation, we sum the infinite series (summation over \( n \)) using Maple. Thus, we arrive at our first destination; the conditional PDF for \( \chi > 0 \) is found in terms of generalized hypergeometric functions.

\[ s^{x+2/3} \exp(-c_\chi s^{2/3}) = \sum_{n=0}^{\infty} \frac{(-c_\chi)^n s^{x+2(n+1)/3}}{n!}. \]  

(38)

where the summation is over the eigenvalues. In Fig. 5, we plot the solution for \( D = 2/5 \) and \( D \rightarrow \infty \) corresponding to the Brownian case. Notice the \( \chi \sim \sqrt{D} \tau^{1/2} \) scaling, which proves the scaling hypothesis Eq. (17). The sum in Eq. (39) converges quickly as long as \( \chi \) is not too large, and so we can use it to construct a plot of \( p(\chi|\tau) \). For example, for Fig. 5, only \( k = 0, \ldots, 4 \) is needed to obtain excellent convergence. We stress that Eq. (39) gives an explicit representation of the scaling function for \( 0 < \chi < \infty \) and hence, by symmetry, for all \( \chi \). The scaling variable is

\[ v_{3/2} = \frac{\chi}{\sqrt{D} \tau^{3/2}}, \]  

(40)

and the scaling function \( B(v_{3/2}) \) in Eq. (17) can be read directly from Eq. (39).
More specifically, the diffusive scaling of
distributions describe areas under constrained Brownian (or Bessel)
$\tau$
path described by Eq. not an excursion. The area under the Brownian meander
this time interval, the motion is described by a meander
function
observables of interest. Having explicit information on
demonstrate the importance of the meander for specific
are symmetric, with mean equal to 0 (since positive and
investigate the packet
FROM THE AREA UNDER THE BESSEL EXCURSION TO …
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A. The Bessel meander

As noted, the last zero crossing of the momentum
process $p(t)$ takes place at a random time $t_n$, and hence
at time $t$, the particle is unlikely to be on the origin of
momentum. Since the particle, by definition, does not cross
the origin in $(t_n, t)$, its momentum remains either positive
or negative in the backward recurrence time interval
$\tau' = t - t_n$ (with equal probability). This means that, in
this time interval, the motion is described by a meander
not an excursion. The area under the Brownian meander
was investigated previously [8,9,57]. For our purposes, we
need to investigate the Bessel meander. This is a Langevin
path described by Eq. (18) constrained to remain positive,
which starts at the origin but is free to end with $p > 0$.
More specifically, the diffusive scaling of $\chi' \sim \tau'^{3/2}$ still
holds. Then, as for the pairs $(\chi, \tau)$, we have the conditional PDF

$$p_M(\chi'|\tau') \sim D^{-1/2}(\tau')^{-3/2}B_M\left(\frac{\chi'}{D(\tau')^{3/2}}\right).$$

Here, the subscript $M$ stands for a meander. The scaling function $B_M(\cdot)$ is different from $B(\cdot)$, though clearly both
are symmetric, with mean equal to 0 (since positive and
negative meanders and excursions are equally probable).
The calculation of $B_M(\cdot)$ runs parallel to that for the
excursion and is presented in Appendix E. Soon we will
demonstrate the importance of the meander for specific
observables of interest. Having explicit information on $B(\cdot)$
and $g(\tau)$, Eq. (11) (see Appendix A for details), and the
scaling form of $B_M(\cdot)$ (see Appendix E), we can now
investigate the packet $P(x, t)$.

VI. MONTROLL-WEISS EQUATION FOR FOURIER-LAPLACE TRANSFORM OF $P(x,t)$

We now use tools developed in the random-walk
community [58–61] to relate the joint PDF of a single
excursion, Eqs. (13) and (14),

$$\psi(\chi, \tau) = g(\tau)D^{-1/2}\tau^{-3/2}B(\chi/(D^{1/2}T^{3/2})), \quad (42)$$

to the probability density $P(x, t)$ for the entire walk. We
find a modified Montroll-Weiss [30,48,62] type of equation
for the Fourier-Laplace transform of $P(x, t)$,

$$\hat{P}(k, u) = \int_{-\infty}^{\infty} dx e^{ikx} \int_{0}^{\infty} d\tau e^{-\mu\tau}P(x, \tau), \quad (43)$$

in terms of the Fourier-Laplace transform of $\psi(\chi, \tau)$, which
will be denoted $\hat{\psi}(k, \mu)$. One modification is that we
include here the correct treatment of the last jump.
Usually, the CTRW model [48] has as an input a single
joint PDF of jump lengths and times, while in our case, we
have essentially two such functions describing the excursions
[i.e. $B(\cdot)$] and the meander [i.e., $B_M(\cdot)$]. Generally,
Montroll-Weiss equations are the starting point for deri-
vation of fractional diffusion equations and the asymptotic
behaviors of the underlying random walks [48]. The
original work of Montroll and Weiss [30] assumed that
there were no correlations between step size $\chi$ and the
waiting time $\tau$, corresponding to a situation called a
decoupled CTRW. The diffusion of atoms in optical lattices
corresponds to a coupled spatial-temporal random-walk
theory first considered by Scher and Lax [63] (see Refs. [64–66]
for recent developments).

We define $\eta_s(x, t)dx$ as the probability that the particle
crossed the momentum state $p = 0$ for the $s$th time in the
time interval $(t, t + dt)$ and that the particle’s position was in
the interval $(x, x + dx)$. This probability is related to the
probability of the previous crossing according to

$$\eta_s(x, t) = \int_{-\infty}^{\infty} dy \int_{0}^{t} d\tau \eta_{s-1}(x - \chi, t - \tau) \times \frac{1}{\sqrt{D\tau^{3/2}}} B\left(\frac{\chi}{\sqrt{D\tau^{3/2}}}\right) g(\tau), \quad (44)$$

where we have used Eq. (42). We change variables
good according to $\chi = v_3/2D^{1/2}T^{3/2}$ and obtain

$$\eta_s(x, t) = \int_{-\infty}^{\infty} dr_{3/2} \int_{0}^{\infty} d\tau \eta_{s-1}(x - v_3/2\sqrt{D\tau^3}, t - \tau) \times B(v_3/2) g(\tau). \quad (45)$$

The process is now described by a sequence of waiting
times $t_1, t_2, \ldots$ and the corresponding generalized velocities
$v_3/2(1), v_3/2(2), \ldots$. The displacement in the $s$th interval is
\[ X_t = v_{3/2}(x) \sqrt{D(t)}^{3/2}. \] (46)

The advantage of this representation of the problem in terms of the pair of microscopic stochastic variables \( \tau, v_{3/2} \) (instead of the correlated pair \( \tau, x \)) is clear from Eq. (45): We may treat \( v_{3/2} \) and \( \tau \) as independent random variables whose corresponding PDFs are \( g(\tau) \) and \( B(v_{3/2}) \), respectively. Here, \( \tau > 0 \) and \(-\infty < v_{3/2} < \infty \). The initial condition \( x = 0 \) at time \( t = 0 \) implies \( \hat{\eta}_0(x, t) = \delta(x) \delta(t) \).

Then, \( P(x, t) \), the probability of finding the particle in \((x, x + dx)\) at time \( t \), is found according to

\[ P(x, t) = \sum_{k=0}^{\infty} \int_{-\infty}^{\infty} \frac{dv_{3/2}}{2\pi} \int_{t}^{\infty} d\tau^* \hat{\eta}_s \left( x - v_{3/2} \sqrt{D(\tau^*)^{3/2}}, t - \tau^* \right) \times B_M(v_{3/2}) W(\tau^*). \] (47)

Here, the survival probability \( W(\tau^*) \equiv 1 - \int_{\tau}^{\infty} g(\tau') d\tau' \) enters since the last jump event took place at \( t - \tau^* \) and in the time period \((t - \tau^*, t)\) the particle did not cross the momentum origin. For this reason, we have in Eq. (47) \( B_M(\cdot) \), not \( B(\cdot) \), since the last time interval in the sequence is a meander and not an excursion. The summation in Eq. (47) is a sum over the number \( s \) of returns to the momentum origin \( p = 0 \). We note that, in this analysis, the particle is always moving, unlike the "wait and then jump" approach used in the original CTRW model.

As usual [48,59], we consider the problem in Laplace-Fourier space where \( t \to u \) and \( x \to k \). Using the convolution theorem and Eq. (45), we find

\[ \hat{\eta}_s(k, u) = \hat{\eta}_{s-1}(k, u) \hat{\mathcal{L}}[B(kD^{1/2} + 3/2)g(\tau)]. \] (48)

Here, \( \hat{\eta}_s(k, u) \) is the Fourier-Laplace transform of \( \eta_s(x, t) \), \( \hat{\mathcal{L}}[\cdot] \) the Fourier transform of \( B(v_{3/2}) \),

\[ \hat{\mathcal{L}}[B(kD^{1/2} + 3/2)g(\tau)] = \int_{-\infty}^{\infty} \exp(-ik\tau^*) B(v_{3/2}) dv_{3/2}. \] (49)

and \( \hat{\mathcal{L}}[\cdot] \) is the Laplace transform operator \( \hat{\mathcal{L}}[h(\tau)] = \int_{0}^{\infty} \exp(-\tau)h(\tau) d\tau \). By definition, the Fourier-Laplace transform of \( \psi(x, \tau) \) is \( \hat{\psi}(k, u) = \hat{\mathcal{L}}[B(kD^{1/2} + 3/2)]g(\tau)] \), and hence

\[ \hat{\eta}_s(k, u) = \hat{\psi}(k, u) \hat{\eta}_{s-1}(k, u). \] (50)

This implies that

\[ \hat{\eta}_s(k, u) = [\hat{\psi}(k, u)]^s, \] (51)

reflecting the renewal property of the underlying random walk. Summing the Fourier-Laplace transform of Eq. (47), applying again the convolution theorem for Fourier and Laplace transforms and using Eq. (51), we find a Montroll-Weiss type of equation for the Fourier-Laplace transform of \( P(x, t) \):

\[ \hat{P}(k, u) = \frac{\hat{\Psi}_M(k, u)}{1 - \hat{\psi}(k, u)}. \] (52)

Here, \( \hat{\Psi}_M(k, u) \) is the Fourier-Laplace transform of \( D^{-1/2}D^{-3/2}B_M(x/\sqrt{D^{3/2}})W(\tau) \). Equation (52) relates statistics of velocity excursions and meanders to the Fourier-Laplace transform of the particle density. The approach is not limited to the specific problem under investigation; namely, the \( x \sim \tau^{1/2} \) scaling is not a necessary condition for the validity of Eq. (52). In the general case, one must revert to \( \psi(x, \tau) = g(\tau)p(\tau) \) instead of the scaling form captured by \( B(\cdot) \). Such an approach might be useful for other systems where the friction is nonlinear.

At this stage, Eq. (52) still depends on \( \epsilon \) since the first passage times for going from \( p = \epsilon \) to the momentum origin, \( p = 0 \), distributed according to \( g(\tau) \), are \( \epsilon \)-dependent (see Appendix A and next section for details). In fact, as \( \epsilon \to 0 \), the number of renewals (i.e., zero crossings) tends to infinity, while in usual CTRWs, the number of renewals (or jumps) is finite for finite observation time \( t \). In the next sections, we will show how the long-time results become independent of \( \epsilon \) in the limit of \( \epsilon \to 0 \).

### VII. ASYMPTOTIC BEHAVIOR OF \( \hat{\psi}(k, u) \)

The Fourier-Laplace transform of the joint distribution of jump times and lengths \( \psi(x, \tau) \) is

\[ \hat{\psi}(k, u) = \int_{-\infty}^{\infty} \frac{d\tau}{\infty} \int_{0}^{\infty} d\tau^* e^{-u\tau + ikr} g(\tau) \frac{1}{\sqrt{D\tau^*}} B\left( \frac{\chi}{\sqrt{D\tau^*}} \right). \] (53)

In this section, we investigate the small \( k \) and small \( u \) behaviors of \( \hat{\psi}(k, u) \), which in the next sections will turn out to be important in the determination of the long-time behavior of \( P(x, t) \). The small \( k \) (\( u \)) limit corresponds to large distance (time), as is well known [48].

#### A. \( k = 0 \) and small \( u \)

Using the normalization condition \( \int_{-\infty}^{\infty} B(v_{3/2}) dv_{3/2} = 1 \), it is easy to verify that

\[ \psi(k, u)|_{k=0} = \hat{g}(u). \] (54)

the Laplace transform of the waiting-time PDF. This waiting-time PDF is investigated analytically in Appendix A. The small \( u \) behavior of \( \hat{g}(u) \) differs depending on the value of \( D \). According to Eq. (11), if \( D < 1 \), the average waiting time \( \langle \tau \rangle \) is finite and so

\[ \hat{g}(u) \sim 1 - \langle \tau \rangle u \quad \text{when} \ D < 1. \] (55)

Here, \( \langle \tau \rangle = \int_{0}^{\infty} \tau g(\tau) d\tau \) is given by a well-known formula for the mean first passage time [50]:

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\[
\langle \tau \rangle = \frac{1}{D} \int_0^\infty dy e^{V(y)/D} \int_y^\infty e^{-V(z)/D} dz, \tag{56}
\]

where \( V(p) \) is the effective momentum potential, Eq. (8).

Equation (56) reflects the absorbing boundary condition at the origin and a reflecting boundary at infinity. Notice that \( \lim_{\tau \to 0} \langle \tau \rangle = 0 \), as is expected, and the leading-order Taylor expansion yields

\[
\langle \tau \rangle \sim \frac{Z}{2D}. \tag{57}
\]

Here, \( Z \) is the normalizing partition function, Eq. (10). Since \( Z \) diverges when \( D \to 1 \) from below, we see that the average waiting time diverges in that limit.

For \( D > 1 \), the mean waiting time is infinite, and so Eq. (55) does not hold and instead, as we show in Appendix A,

\[
\hat{g}(\alpha) \sim 1 - g^* \Gamma(-\alpha)|\alpha|^\alpha, \quad 1/2 < \alpha < 1, \tag{58}
\]

with

\[
g^* = \frac{2\alpha}{(2D)^{3\alpha} \Gamma(3\alpha)} \epsilon. \tag{59}
\]

For large \( \tau \), this yields the power-law behavior in Eq. (11).

\( \hat{g}(\alpha) \sim g^* \Gamma(-\alpha)|\alpha|^\alpha \), which in fact also describes the tail for \( \alpha > 1 \).

The prefactor \( g^* \) vanishes as \( \epsilon \to 0 \) and, importantly, does not depend on the full shape of the effective potential \( V(p) \) but rather only on the value of the dimensionless parameter \( D \). The case \( D = 1 \) contains logarithmic corrections and will not be discussed here. Using Eqs. (58) and (59), we find nontrivial distributions of the time-interval straddling time \( t \) and the backward and forward recurrence times, which were previously treated by mathematicians without the \( \epsilon \) trick. This connection between renewal theory and the statistics of zero-crossings of the Bessel process is discussed in Appendix F.

**B. The \( u = 0 \), small \( k \) limit**

Similarly to Eq. (54), by definition,

\[
\hat{\psi}(k,u)|_{u=0} = \hat{q}(k), \tag{60}
\]

where \( \hat{q}(k) \) is the Fourier transform of the symmetric jump length distribution \( q(x) \). Alternatively, we can use Eq. (53) with \( u = 0 \) and then, upon changing variables according to \( x/(\sqrt{D}e^{3/2}) = v_{3/2} \) and using the normalization of the scaling function \( B(v_{3/2}) \), we find

\[
\hat{\psi}(k,u)|_{u=0} = 1 - \int_0^\infty dr \int_{v_{3/2}}^\infty dv_{3/2} \left( 1 - e^{i v_{3/2} \sqrt{D} \omega} \right) \hat{g}(\omega) \times B'(v_{3/2}). \tag{61}
\]

Using \( \hat{B}(k) \), the Fourier transform of \( B(v_{3/2}) \), Eq. (49), we can rewrite Eq. (61) as

\[
\hat{\psi}(k,u)|_{u=0} = 1 - \int_0^\infty dr [1 - \hat{B}(k\sqrt{D} e^{3/2})] \hat{g}(\omega). \tag{62}
\]

This expression is the starting point for the small \( k \) expansion carried out in Appendix C, which gives for \( \nu < 2, \ (D > 1/5) \),

\[
\hat{\psi}(k,u)|_{u=0} \sim 1 - g^*_\nu \left( \nu \frac{\nu}{2} \right)^\nu \cos \left( \frac{\nu \pi}{2} \right) \sqrt{D} \nu. \tag{63}
\]

The nonanalytical character of this expansion is responsible for the anomalous diffusion of the atom’s position. The first term on the right-hand side is the normalization condition, and the second, proportional to \( \sqrt{D} \nu \), is consistent with the fat-tailed PDF of jump lengths \( q(x) \alpha x^{-1+\nu} \), Eq. (11), namely, an excursion length whose variance diverges since \( \nu < 2 \). As expected, this behavior is in full agreement with Eq. (11) since \( 4/3 + \beta = 1 + \nu \).

In Eq. (63), we see the noninteger moment of the scaling function \( B(\cdot) \),

\[
\langle |v_{3/2}|^\nu \rangle = \int_{-\infty}^{\infty} |v_{3/2}|^\nu B(v_{3/2}) dv_{3/2}. \tag{64}
\]

Given the formidable structure of the scaling function \( B(v_{3/2}) \), we do not describe here [67] the direct method to obtain noninteger moments like Eq. (64). Instead, we present a method that gives \( \langle |v_{3/2}|^\nu \rangle \) indirectly. In a future publication [67], we will discuss this and other moments of \( B(v_{3/2}) \).

We can use Eq. (60) together with \( q(x) \sim q_\nu x^{-1+\nu} \) to find, in Fourier space,

\[
\hat{q}(k) \sim 1 - \frac{\pi q_\nu |k|^\nu}{\sin(\frac{\pi}{2} \Gamma(1+\nu))}. \tag{65}
\]

for \( \nu < 2 \). In Appendix B, we investigate the area under a Bessel excursion regularized at the origin using a backward Fokker-Planck equation [6], which gives the amplitude of jump lengths

\[
q_\nu \sim \epsilon \frac{\nu}{3\nu - 1} \left( \frac{3
u - 1}{2\Gamma(\nu)} \right). \tag{66}
\]

Comparing with Eq. (63), we arrive at the following simple relation:

\[
\langle |v_{3/2}|^\nu \rangle = \lim_{\epsilon \to 0} \frac{3 q_\nu}{D^{3/2} \epsilon}. \tag{67}
\]

Thus, we may use Eqs. (59) and (66) to find the desired \( \epsilon \)th moment of the scaling function \( B(\cdot) \),

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\[ \langle |v_{3/2}|^p \rangle = \frac{2^{2-p-1} \Gamma(\alpha)}{3^{2-p} \Gamma(\nu)}. \] \hspace{1cm} (68)

In the limit \( D \to \infty \), we have \( \nu = 1/3 \), and then

\[ \lim_{D \to \infty} \langle |v_{3/2}|^p \rangle = \frac{\sqrt{3}^{4/3}}{2 \Gamma(1/3)} = 1.431. \] \hspace{1cm} (69)

We see that while the amplitudes \( q^* \) and \( q^* \) vanish as the convenient theoretical tool \( e \to 0 \), \( \langle |v_{3/2}|^p \rangle \) is independent of it in the limit, indicating the usefulness of this variable.

**C. The second moment \( \langle (v_{3/2})^2 \rangle \)**

As we show below, the second moment \( \langle (v_{3/2})^2 \rangle \) determines the mean-square displacement of the particles. The mean vanishes since, in the underlying random walk, positive and negative excursions are equally likely. As for the \( r^{th} \) moment, \( \langle |v_{3/2}|^p \rangle \), the extraction of integer moments from the exact solution is not straightforward. In the regime \( 1/5 < D < \infty \), an excellent approximation is

\[ \langle (v_{3/2})^2 \rangle = \frac{5}{6} + \frac{83}{270D}. \] \hspace{1cm} (70)

Similarly, for the meander, we find

\[ \langle (v_{3/2})^2 \rangle_M = \frac{1}{3D} + \frac{59}{30}. \] \hspace{1cm} (71)

This and Eq. (70) agree with known results in the Brownian limit \( D \to \infty \) [8] (units and notations used in Ref. [8] are not used by us). In Fig. 6, we show that \( \langle (v_{3/2})^2 \rangle \) and \( \langle (v_{3/2})^2 \rangle_M \) nicely match their linear approximations [67].

**VIII. THE LÉVY PHASE**

We now explain why Lévy statistics, and hence the generalized central-limit theorem, describes the central part of the diffusion profile \( P(x, t) \) at long times for \( 1/5 < D < 1 \). The Lévy profile is cut off in the tails of the distribution because of the correlations between jump length and time investigated here. We focus on the Lévy phase first because it has been reported experimentally [3].

The key idea is that in the regime \( D < 1 \), \( \langle \tau \rangle \) is finite, and hence the number of jumps \( n \) scales with \( t/\langle \tau \rangle \) when \( t \) is large [47]. At the same time, the jump length PDF still does not have a variance (since \( 1/5 < D \) ), which means that the usual Gaussian central-limit theorem does not hold. Instead, because of the power-law distribution \( q(x) \propto |x|^{-(1+\nu)} \), the process belongs to the domain of attraction of a Lévy stable law. As long as \( x \) is not too large, the correlations are not important. However, when \( x \propto t^{3/2} \), the simple Lévy picture breaks down since clearly we cannot perform a jump larger than the order of \( t^{3/2} \). Thus, for \( 1/5 < D < 1 \) and \( t^{3/2} \ll k \ll 1 \), we can approximate

\[ \hat{\psi}(k, u) = 1 - u\langle \tau \rangle + c_0 q_x |k|^\nu \ldots, \] \hspace{1cm} (72)

where we have used Eqs. (55) and (65). Here, from Eq. (65) \( c_0 = \pi [\sin(\pi \nu/2)\Gamma(1 + \nu)] \), Equation (72) corresponds to a decoupling scheme, \( \hat{\psi}(k, u) = \hat{\eta}(u)\hat{Q}(k) \), which, according to arguments in Ref. [60], is exact in the long-time limit in the regime under investigation. Notice that \( 1/5 < D < 1 \) gives \( 2/3 < \nu \ll 2 \).

Using the Montroll-Weiss-type equation (52) and \( \hat{u}_M(k, u) \sim \langle \tau \rangle \), which is easy to prove, we find the Fourier-Laplace representation of the solution

\[ \hat{P}(k, u) \sim \frac{\langle \tau \rangle}{u\langle \tau \rangle + c_0 q_x |k|^\nu}. \] \hspace{1cm} (73)

As mentioned, both \( \langle \tau \rangle \) and \( q_x \) vanish as \( e \) approaches zero. Rearranging, we have

\[ \hat{P}(k, u) \sim \frac{1}{u + K_e |k|^\nu}, \] \hspace{1cm} (74)

where

\[ K_e = \lim_{e \to 0} \frac{c_0 q_x}{\langle \tau \rangle}, \] \hspace{1cm} (75)

and from Eqs. (57) and (66) [26],

\[ K_e = \lim_{e \to 0} \frac{c_0 q_x}{\langle \tau \rangle} = \frac{1}{\nu} \pi \frac{3\nu - 1}{3^{2-\nu} \Gamma(\nu + 1)} |k|^\nu. \] \hspace{1cm} (76)

\( K_e \) is called the anomalous diffusion coefficient. When returning to physical units, we get

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which has units cm$^2$/sec. An equivalent expression is $K_x = c_0(|x^2/2|)^{3/2} \text{lim}_{\alpha \to 0^+} e^{3/2}$. \( \tilde{P}(k,u) \) as given in Eq. (74) is, in fact, precisely the symmetric Lévy PDF in Laplace-Fourier space, whose \((x,t)\) presentation [see Eq. (B17) of [62]] is

$$P(x,t) \sim \frac{1}{(K_x t)^{1/2}} L_{\nu,0} \left[ \frac{x}{(K_x t)^{1/2}} \right]$$  \( (78) \)

for $2/3 < \nu < 2$. The properties of the Lévy function $L_{\nu,0}(\bullet)$ are well known. The Fourier transform of this solution is $\text{exp}(-K_x |k|^2)$ for $2/3 < \nu < 2$, which can serve as the working definition of the solution via the inverse Fourier transform.

Figure 7 shows excellent agreement between simulations and the theory. It also illustrates the cutoff on Lévy statistics, which is found at distances $x \propto \tau^{1/2}$. Beyond this length scale, the density falls off rapidly. This, as noted above, is the result of the correlation between $\chi$ and $\tau$, as there is essentially no weight associated with paths whose displacement is greater than the order of $\tau^{1/2}$.

This cutoff ensures the finiteness of the mean-square displacement (see Fig. 8). Using the power-law tail of the Lévy PDF $L_{\nu}(x) \propto x^{-(1+\nu)}$ and the time scaling of the cutoff, we get

$$\langle x^2 \rangle = \int_0^{\tau^{2/3}} t^{-(1/\nu)}(x/t^{1/\nu})^{-(1+\nu)} x^2 dx \propto t^{\nu-3\nu/2}$$  \( (79) \)

for $2/3 < \nu < 2 (1/5 < D < 1)$. If we were to rely only on the Lévy PDF, Eq. (78), we would get $\langle x^2 \rangle = \infty$. Thus, the Lévy PDF solution must be treated with care, realizing its limitations in the statistical description of the moments of the distribution and its tails. As we soon show, if $\nu > 2$, we get normal diffusion $\langle x^2 \rangle \propto t$, while for $\nu < 2/3$, we have $\langle x^2 \rangle \propto t^{\nu}$. This last behavior is due to the correlations, which restrict jumps longer than $t^{1/2}$. So, we have three phases of motion [26]:

$$\langle x^2 \rangle \propto \begin{cases} t & 2 < \nu \\ t^{\nu-3\nu/2} & 2/3 < \nu < 2 \\ t^3 & 1/3 < \nu < 2/3 \end{cases} \text{ Obukhov-Richardson.}$$  \( (80) \)

These simple scaling arguments for the mean-square displacement can be derived from a more rigorous first-principle approach, to which we will turn in Sec. IX.

A. The diffusion exponent

Wickenbrock et al. [2] investigated the additional effect of a high-frequency (HF) oscillating force $F_{HF} = A_{HF} \sin(\omega_{HF}t + \phi_0)$ on the dynamics of the atoms, where the frequency $\omega_{HF}$ is much larger than other frequencies in the system. According to Ref. [2], in the limit of a strong drive, the depth of the optical lattice potential is renormalized, $U_0 \to U_0 |J_0(2kr)|$, where $J_0(\cdot)$ is the Bessel function of the first kind, $r = A_{HF}/(m a_{HF}^2)$, and $k$ is the laser-field wave vector. This elegant setup allows the control of the transport via the renormalization of the optical depth $U_0$ and, according to Eq. (7), the control of the dimensionless parameter $D$. For example, in the vicinity of the zeros of the Bessel function $J_0$, we clearly find an effective shallow lattice, which, according to the theory, corresponds to the
we find losses become insignificant. As mentioned in our summary, we found in the previous section since 

\[ P(k, u) \sim \frac{1 - \frac{1}{2} D \langle (\nu_{3/2})^2 \rangle_u f_1(u) k^2}{u} + \frac{1}{2} D \langle (\nu_{3/2})^2 \rangle_u f_2(u) k^2. \]

This second-order expansion contains the second-order moment of the scaling function \( \langle (\nu_{3/2})^2 \rangle_u \) and similarly for \( \langle (\nu_{3/2})^2 \rangle_u \). Thus, the numerator (denominator) in Eq. (81) contains a term describing the meander (excursions) contribution. Hence, from symmetry \( \langle \nu_{3/2} \rangle = 0 \), the expansion does not contain linear terms. Here,

\[ f_1(u) = -\frac{1}{\hat{W}(u)} \frac{d^3}{du^3} \hat{W}(u), \]

where \( \hat{W}(u) = [1 - \hat{g}(u)]/u \) is the Laplace transform of the survival probability \( W(\tau) \) and

\[ f_2(u) = -\frac{d}{du} \frac{\hat{g}(u)}{1 - \hat{g}(u)}. \]

The third-order derivative with respect to \( u \) is clearly related to the \( \chi \propto r^{3/2} \) scaling we have found and to the second-order expansion. While the \( k^2 \) expansion in Eq. (81) works fine for small \( k \) and finite \( u \), when \( u \to 0 \) we get divergences. For example, when \( \alpha < 1 \) we have \( \hat{g}(u) \sim 1 - \hat{g}(u) \), and hence the third-order derivative of \( \hat{g}(u) \) diverges as \( u^{-3} \) as \( u \to 0 \). For \( 1 < \alpha < 3 \), \( f_2(u) \) diverges as \( u^{\alpha-4} \) and for \( \alpha > 3 \), \( f_2(u) \) diverges as \( u^{\alpha-1} \). This change in scaling behavior at \( \alpha = 3 \), as we shall presently see, marks the transition from anomalous diffusion to normal, consistent with what we found in the previous section since \( \alpha = 3 \) gives \( D = 1/5 \) and hence \( \nu = 2 \). Of course, the \( k^2 \) behavior in Eq. (81) is very different from the nonanalytical \( |k|^\alpha \) found in Eq. (63). This indicates that the order of taking the limits \( k \to 0 \) and \( u \to 0 \) is noncommuting [59].

The Laplace transform of the mean-square displacement of the atoms is given by

\[ \langle \hat{x}^2(u) \rangle = -\frac{d^2}{du^2} \langle \hat{P}(k, u) \rangle \bigg|_{k=0}. \]

Hence, for particles starting on the origin, one can easily see that

\[ \hat{x}^2(u) \propto u^{-\omega/2} \]

FIG. 9. The anomalous diffusion exponent \( \langle \hat{x}^2 \rangle \sim r^{2} \) versus the high-frequency control parameter \( kr \) defined in the text. The triangles and diamonds are simulation results from Fig. 1 in Ref. [2] with \( U_0 = 200E_c \) (with varying photon scattering rate; see details in Ref. [2]). The curve is the theory, Eqs. (7) and (80), with a renormalized \( U_0 \) as explained in the text with \( c = 34.7 \). The transition between normal \( 2 \xi = 1 \), Lévy \( 1 < 2 \xi < 3 \), and Richardson \( 2 \xi = 3 \) behaviors is clearly visible, and as shown by Wickenbrock et al., it is controlled by the high-frequency field.

Richardson phase. In the experiment [2], resonances in the transport are observed close to the zeros of the Bessel functions, namely, an enhanced spreading of the atoms. However, as pointed out in Ref. [2], many superdiffusing atoms are lost, which leads to an underestimate of the diffusion exponent. In Ref. [2], the diffusion exponent was found using Monte Carlo simulation (see Fig. 1 of that reference). To demonstrate the predictive power of our theory, at least for exponents, we compare between theory and the numerics [2]. As mentioned in our summary, we postpone a comparison of experiments to theory until the losses become insignificant.

From Eq. (80) and within the Lévy phase, we find superdiffusion and \( \langle \hat{x}^2 \rangle \propto r^{2} \) with \( 2 \xi = 4 - 3\nu/2 = 7/2 - 1/(2D) \) for \( 1/5 < D < 1 \) so \( 1 < 2 \xi < 3 \). For \( D > 1 \), we find \( \langle \hat{x}^2 \rangle \propto r^{3} \), so \( 2 \xi = 3 \) and normal diffusion with \( 2 \xi = 1 \) when \( D < 1/5 \). With a renormalized Eq. (7), \( D = cE_R/[U_0|J_0(2kr)|] \), we present in Fig. 9 the mean-square displacement exponent versus the control parameter \( kr \), with nice agreement between theory and simulations.

IX. THE MEAN-SQUARE DISPLACEMENT

Here, we present the calculation of the mean-square displacement of the atoms using the Montroll-Weiss equation. Our aim, as declared in the Introduction, is to unravel the quantitative connections between the transport and the statistics of excursions, for example, the relation between the mean-square displacement and moments of the area under the Bessel excursion and meander. Such relations are expected to be general beyond the model under investigation. A different strategy for the calculation can be based on a super-aging velocity-velocity correlation function approach [21,36].

To derive Eq. (78), we assumed in Eq. (74) that \( u \) and \( K_\nu[|k|^\alpha] \) are of the same order of magnitude. We now consider a different small \( k \), \( u \) limit of \( \hat{P}(k, u) \). We first expand the numerator and denominator of \( \hat{P}(k, u) \) in the small parameter \( k \) to second order. We leave \( u \) fixed and find

\[ \hat{P}(k, u) \sim \frac{1 - \frac{1}{2} D \langle (\nu_{3/2})^2 \rangle_u f_1(u) k^2}{u + \frac{1}{2} D \langle (\nu_{3/2})^2 \rangle_u f_2(u) k^2}. \]
\[ \langle \hat{x}^2(u) \rangle = \lim_{\epsilon \to 0} \frac{D}{u} \left[ \langle (v_{3/2})^2 \rangle_M \hat{f}_1(u) + \langle (v_{3/2})^2 \rangle_M \hat{f}_2(u) \right]. \]  

(85)

The second moment is finite because of the observed fast decay of the scaling function \( B(v_{3/2}) \) for \( v_{3/2} \gg 1 \), ensuring that \( \langle (v_{3/2})^2 \rangle \) and \( \langle (v_{3/2})^2 \rangle_M \) are both finite.

### A. Obukhov-Richardson diffusion

When \( \alpha < 1 \) (\( D > 1 \)), we have \( \hat{g}(u) \sim 1 - Gu^\alpha + \cdots \), where \( G = g^\prime \Gamma(-\alpha) \). Using Eqs. (82) and (83), \( \hat{f}_1(u) \sim c_1 u^{-3} \) and \( \hat{f}_2 \sim c_2 u^{-3} \) for small \( u \) with

\[
\begin{align*}
    c_1 &= (1 - \alpha)(2 - \alpha)(3 - \alpha), \\
    c_2 &= \alpha(1 - \alpha)(2 - \alpha).
\end{align*}
\]

(86)

The small \( k, u \) expansion of \( \hat{P}(k, u) \), Eq. (81), is

\[ \hat{P}(k, u) \sim \frac{1 - \frac{1}{2} D \langle (v_{3/2})^2 \rangle_M c_1 u^{-3} k^2}{u + \frac{1}{2} D \langle (v_{3/2})^2 \rangle c_2 u^{-3} k^2}. \]

(87)

which is \( g^\prime \) and \( \epsilon \) independent. The mean-square displacement in the small \( u \) limit is

\[ \langle \hat{x}^2(u) \rangle \sim D \langle c_1 \langle (v_{3/2})^2 \rangle_M + c_2 \langle (v_{3/2})^2 \rangle \rangle_\epsilon u^{-4}. \]

(88)

Converting to the time domain,

\[ \langle x^2(t) \rangle \sim \left[ c_1 \langle (v_{3/2})^2 \rangle_M + c_2 \langle (v_{3/2})^2 \rangle \right] \frac{D t^3}{6}. \]

(89)

The scaling \( x^2 \propto t^3 \) in this \( \alpha < 1 \) regime is similar to Richardson’s observations concerning the relative diffusion of a pair of particles in turbulence (see discussion below). We see that in this regime, both the meander and the excursion contribute to the computation of the mean-square displacement. The theory agrees with the finite time simulations presented in Fig. 10, where we see that \( \langle x^2 \rangle / t^2 \) approaches zero as \( D \to 1 \) in the long-time limit.

We can verify Eq. (89) in the limit of large \( D \), where the friction force is negligible. In that case, Eq. (4) with \( F(p) = 0 \) easily gives \( \langle x^2 \rangle = 2 D t^3 / 3 \). On the other hand, we have, for Brownian excursions and meanders, \( \langle (v_{3/2})^2 \rangle = 5 / 6 \) and \( \langle (v_{3/2})^2 \rangle_M = 59 / 30 \), Eqs. (70) and (71) [8]. Using \( \lim_{D \to \infty} \alpha = 1 / 2 \), we have, in this limit, \( c_1 = 15 / 8 \) and \( c_2 = 3 / 8 \). Plugging these numbers into Eq. (89), we get \( \langle x^2 \rangle = 2 D t^3 / 3 \), as is expected. This simple demonstration implies that it is essential to treat the last jump event properly (as a meander), and previous CTRW approaches relying on a unique jump length distribution lead to wrong conclusions [i.e., replacing \( \langle (v_{3/2})^2 \rangle_M \) with \( \langle (v_{3/2})^2 \rangle \) is wrong]. Furthermore, in this limit, the contribution of the meander is numerically larger than the contribution of the excursions, even though the number of excursions is large. Notice, however, that for the calculation of the Lévy density \( P(x, t) \), Eq. (78), the statistics of the meander do not enter. Hence, depending on the observable of interest, and the value of \( D \), the meander may be either a relevant part of the theory or not.

### B. Superdiffusion

For \( 1 < \alpha < 2 \) (\( 1 / 3 < D < 1 \)), we use the expansion

\[ \hat{g}(u) = 1 - u(\tau) + Gu^\alpha + \cdots. \]

(90)

Here, the first moment of the waiting time is finite, while the second moment diverges. Crucially, both \( \langle \hat{r} \rangle \) and \( G = g^\prime \Gamma(-\alpha) \) vanish as \( \epsilon \to 0 \). For this parameter regime, we have found Lévy behavior for the central part of \( P(x, t) \). Using Eq. (85), we find

\[ \langle \hat{x}^2(u) \rangle \sim D \lim_{\epsilon \to 0} \left[ c_1 \langle (v_{3/2})^2 \rangle_M + c_2 \langle (v_{3/2})^2 \rangle \right] \frac{G}{\langle \tau \rangle} u^{\alpha - 5}. \]

(91)

Using Eqs. (57) and (59),

\[ \lim_{\epsilon \to 0} \frac{G}{\langle \tau \rangle} = \frac{2D}{\Gamma(2\alpha)} \frac{2\alpha}{\Gamma(\alpha)} \]

(92)

Inverting Eq. (91) to the time domain and inserting \( D = (2\alpha - 1)^{-1} \) gives

\[ \langle x^2(t) \rangle \sim \left[ \alpha D \Gamma(\alpha) \right] \frac{2D}{\Gamma(2\alpha)} \left[ c_1 \langle (v_{3/2})^2 \rangle_M + c_2 \langle (v_{3/2})^2 \rangle \right] t^{\alpha - 5}. \]

(93)
The same result is valid for $2 < \alpha < 3$. This behavior depends on the normalizing partition function $Z$; namely, the shape of the potential $V(p)$ in the vicinity of the momentum origin becomes important, unlike the Richardson-Obukhov phase, where only the large $p$ behavior of $V(p)$ is important in the long-time limit [i.e., the case $\alpha < 1$, Eq. (89)]. Notice that $|c_1|$ in Eq. (93) tends to zero when $\alpha \to 3$ from below. This means that the meander becomes irrelevant when we approach the normal diffusion phase $x^2 \sim t$. Figure 11 compares simulation and theory and demonstrates that $(\chi^2)/t^{5/3}$ diverges as the transition to the Gaussian phase $D < 1/5$ is approached.

C. Breakdown of scaling assumption—normal diffusion

Our starting point was the scaling hypothesis $\chi^2 \sim t^3$, Eq. (14) and Fig. 3. Indeed, we have shown that for large $\chi$ and $\tau$, the scaling function $B(.)$ describes the conditional probability density $p(\chi|\tau)$ for all values of $D$. However, this does not imply that the scaling solution $B(.)$ is always relevant for the calculation of the particle density $P(x,t)$. Roughly speaking, so far we have assumed that large jumps $\chi$ and long waiting times $\tau$ dominate the underlying process. No one guarantees, however, that this large $\chi$ limit (long jumps) is important while small $\chi$ (small jumps) can be neglected. Indeed, when we switch over to the normal diffusion phase $D < 1/5$, the small-scale aspects of the process become important [meaning that the regularization of the potential $V(p)$ when $p \to 0$ becomes crucial]. This is similar to the Lévy versus Gaussian central-limit theorem arguments, where the former is controlled by the tails of the distribution while the latter is controlled by the variance. Let us see this breakdown of scaling in more detail.

When $\alpha > 3$, we have the expansion $\hat{g}(u) \sim 1 - u(\tau) + u^2(\tau^3)/2 - u^3(\tau^5)/6 + \cdots$, where the first three integer moments of the waiting-time PDF are finite (see Appendix A). Then, the function $\hat{f}_1(u)$ is negligible when $u \to 0$, while $\hat{f}_2(u) \sim (\tau^3)/u(\tau)$. We find

$$\langle \chi^2 \rangle \sim D\langle (r^{3/2})^2 \rangle \lim_{\epsilon \to 0} \frac{\langle \tau^3 \rangle}{\langle \tau \rangle} t. \quad (94)$$

Thus, as stated above, $\alpha = 3$ (or $D = 1/5$) marks the transition between the anomalous superdiffusive phase and the normal diffusion phase. Furthermore, in this case, the meander is of no importance.

However, Eq. (94) is correct only if our assumptions about scaling are valid. By definition,

$$\langle \chi^2 \rangle = \int_0^\infty \int_0^\infty \psi(\chi, \tau) \chi^2 d\chi d\tau, \quad (95)$$

and if the scaling hypothesis holds,

$$\langle \chi^2 \rangle = \int_0^\infty \int_{-\infty}^{\infty} g(\tau) \frac{B(\chi/\sqrt{D\tau^{3/2}})}{\sqrt{D\tau^{3/2}}} \chi^2 d\chi d\tau, \quad (96)$$

which gives

$$\langle \chi^2 \rangle = D\langle (r^{3/2})^2 \rangle \langle \tau \rangle. \quad (97)$$

Inserting Eq. (97) into (94), we get the expected result discussed in the Introduction,

$$\langle \chi^2 \rangle \sim \lim_{\epsilon \to 0} \frac{\langle \chi^2 \rangle}{\langle \tau \rangle} t. \quad (98)$$

This simple result is the correct one, even though the scaling assumption is not valid in this regime. Sometimes
we reach truthful conclusions, even though the assumptions along the way are invalid. To see the breakdown of scaling, we plot, in Fig. 12, $Y_{\text{dev}} = \lim_{t \to 0} \langle \chi(t) \rangle = [D(x_0^2)]^{1/2}$ versus $1/D$, which should be zero in the normal phase if the scaling hypothesis holds (where $\langle \chi^2 \rangle$ and $\langle \tau \rangle$ are finite). The calculations of $\langle \chi^2 \rangle$ and $\langle \tau \rangle$ are given in Appendices A and B. We see that at the transition point, $D = 1/5$, $Y_{\text{dev}} = 0$, but otherwise $Y_{\text{dev}} \neq 0$. Hence, the scaling hypothesis does not work in the normal phase $D < 1/5$. This means we need another approach for normal diffusion, which, luckily, is easy to handle. For $D > 1/5$, long jumps dominate since $\langle \chi^2 \rangle$ diverges, and then our scaling theory works fine.

X. THE RICHARDSON-OBUKHOV PHASE

In Sec. VIII, we discussed the Lévy phase that is found for $1/5 < D < 1$. When the average jump duration $\langle \tau \rangle$ diverges, i.e., for $D > 1$, the dynamics of $P(x, t)$ enters a new phase. Since the index of the Lévy PDF $\nu$ approaches $2/3$ as $D$ approaches 1, $x$ scales like $r^{1/2}$ in the limit. Because of the correlations between $x$ and $r$, $x$ cannot grow faster than this, so in this regime, $P(x, t) \sim r^{-3/2} \delta(x/r^{1/2})$. An example for this behavior is presented in Fig. 13. This scaling is that of free diffusion; namely, momentum scales like $p \sim t^{1/2}$, and hence the time integral over the momentum scales like $x \sim t^{3/2}$. Indeed, in the absence of the logarithmic potential, namely, in the limit $D \gg 1$, the Langevin equations (4) give

$$P(x, t) \sim \sqrt{\frac{3}{4\pi D t}} \exp \left( -\frac{3x^2}{4Dt} \right). \quad (99)$$

This limit describes the Obukhov model for a tracer particle in turbulent flow, where the velocity follows simple Brownian motion [68,69]. The scalings are related to Kolmogorov’s theory of 1941 [see Eq. (3) in Ref. [69]] and to Richardson’s diffusion, $\langle x^2 \rangle \sim t^\nu$ [70]. Equation (99) is valid when the optical potential depth is small since $D \to \infty$ when $U_0 \to 0$. This limit should be taken with care, as the observation time must be made large before considering the limit of weak potential. In the opposite scenario, $U_0 \to 0$ before $t \to \infty$, we expect ballistic motion $|x| \sim \tau$ since in this case the optical lattice has not had time to make itself felt [3]. Physically, the atoms in this phase perform a random walk in momentum space due to random emission events, which in turn give the $x \sim t^{1/2}$ scaling. For shallow lattices, the Sisyphus cooling mechanism breaks down, in the sense that transitions from maximum to minimum of the potential field created by the laser fields are not preferred over the inverse transitions. Thus, the deterministic dissipation is not effective, and we are left with Brownian scaling in momentum space, $p \sim t^{1/2}$.

XI. THE NORMAL PHASE

When the variance of the jump length is finite, namely, $\nu > 2$ ($D < 1/5$), we get normal diffusion. Here, the variance of jump lengths is finite, and hence the scale-free dynamics breaks down. The breakdown of scaling means that instead of using the scaling function $B(\cdot)$, e.g., in Eq. (42) for $\psi(x, t)$, we must use the joint PDF $\psi(x, t) = g(t)p(x|\tau)$, Eq. (13), and, in principle, not limit ourselves to large $\tau$. However, luckily there is no need for a new calculation.

We focus on the central part of the density $P(x, t)$, where central-limit theorem arguments hold. In this normal case, the spatiotemporal distribution of jump times and jump lengths effectively decouples, similar to the Lévy phase. Since the variance of jump size and the averaged time for jumps are finite, many small jumps contribute to the total displacement, and hence in the long-time limit, we expect Gaussian behavior with no correlations between jump lengths and waiting times; i.e., the decoupling approximation is expected to work. More precisely, the average waiting time is finite, so $\bar{\chi}(u) \sim 1 - \langle \chi \rangle u + \cdots$; the variance of jump lengths is also finite, so the Fourier transform of $q(\xi)$ has the following small $k$ expansion: $\tilde{q}(k) = 1 - \langle \chi^2 \rangle k^2/2 + \cdots$, where $\langle \chi^2 \rangle = \int_0^\infty x^2 q(x) dx$ is the variance of the jump lengths. This variance is investigated in Appendix B using a backward Fokker-Planck equation. In the small $k, u$ limit, $\tilde{\psi}(k, u) \sim 1 - u (\xi + k^2/2 + \cdots)$, and the Montroll-Weiss equation (52) is

$$\tilde{P}(k, u) \sim \frac{1}{u + K_1 k^2}.$$ \quad (100)

This is the expected Gaussian behavior for the position probability density.
\[ P(x, t) \sim \frac{1}{\sqrt{4\pi K_2 t}} \exp \left( -\frac{x^2}{4K_2 t} \right) \]  
(101)

and

\[ K_2 = \lim_{\epsilon \to 0} \langle \chi^2 \rangle \]  
(102)

is the diffusion constant, namely,

\[ \langle \chi^2(t) \rangle \sim 2K_2 t. \]  
(103)

Equation (102) relates the statistics of the excursions, i.e., the variance of the area under the excursion \( \langle \chi^2 \rangle \) and its average duration \( \langle \tau \rangle \) to the diffusion constant (here, \( \langle \chi \rangle = 0 \)). As noted in the Introduction, the equation has the structure of the famous Einstein relation, relating the variance of jump size and the time between jumps to the diffusion.

While \( \langle \tau \rangle \) [Eq. (57)] and \( \langle \chi^2 \rangle \) [Eq. (B22)] approach zero when \( \epsilon \to 0 \), their ratio remains finite and gives

\[ K_2 = \frac{1}{D} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \exp \left( -\frac{p^2}{D} \right) \left[ \int_{0}^{\infty} dp' e^{-V(p')/D} p' \right]^2. \]  
(104)

This equation was derived previously using different approaches [20,36]. The diffusion constant \( K_2 \) diverges as \( D \to 1/5 \) from below, indicating the transition to the superdiffusive phase. A sharp increase in the diffusion constant \( K_2 \) as the intensity of the laser reaches a critical value was demonstrated experimentally [20]. Equation (104) can be derived using the Green-Kubo formalism [6], so in the normal phase, the analysis of the statistics of excursions is an alternative to the usual methods. It seems that in the Lévy phase, the analysis of statistics of excursions is vital. Specifically, the usual Green-Kubo formalism breaks down since \( K_2 \) is infinite, and the calculation of the anomalous diffusion coefficient \( K_2 \) cannot be based on the computation of a stationary velocity correlation function.

XII. DISCUSSION

Starting with the Langevin description of the semiclassical motion of atoms undergoing a Sisyphus cooling process, we mapped the problem onto a random-walk scheme using excursions as a tool. Thus, our work combines ideas from the theory of stochastic processes with cold atom physics. We now summarize the key ingredients of our results and their predictions from the point of view of these two communities.

A. Fractional diffusion equation and CTRW theory revisited

The first ingredient of the theory was the calculation of the coupled joint distribution \( \psi(\chi, \tau) \). Rather generally, the coupled joint distribution of waiting-time and jump length is the microscopic ingredient for a coupled continuous-time random walk [48,58–63]. These distributions, however, are difficult to obtain from first-principle calculations and are usually treated in a simplified manner. For example, postulating \( \psi(\chi, \tau) = g(\tau)\delta(\chi - \tau)/2 \) is common in stochastic theories. Hence, we provide an important pillar for the foundation of this widely applied approach.

Our work gives the relation between velocity excursions and random-walk theory and hence diffusion phenomena. Since zero crossing in velocity space is obviously a very general feature of physical paths of random processes, we expect that the areal distributions of excursions and meanders will play an important role in other systems, at least as a tool for the calculation of transport and diffusion (e.g., in principle, our approach can also treat the case of a constant applied force, where a mean net flow is induced). The celebrated Montroll-Weiss equation needs two important modifications. First, since the underlying process is continuous, we regularized the process, such that the excursions and meander start at \( \epsilon \to 0 \). Second, the statistics of the last jump event, i.e., the meander, must be treated with care. In contrast, usually it is assumed that only \( \psi(\chi, \tau) \) is needed for a microscopical description of a continuous-time random-walk model. This and the correlations between \( \chi \) and \( \tau \) make the problem challenging, and usual approaches to diffusion fail. With our approach, the behavior of the packet is mapped onto a problem of the calculation of areas under Bessel excursions and meanders. For example, we derived the relation between the mean-square displacement of the atoms and the areas under both the Bessel excursion and meander, Eqs. (85), (89), and (93). A decoupled scheme that neglects the correlations \( \psi(\chi, \tau) = g(\tau)q(\chi) \) gives a diverging result for \( \nu < 2 \), which is unphysical.

In the regime \( 1/5 < D < 1 \), which we have called the Lévy regime, the correlations between jump lengths and waiting times point to the limitations of the fractional diffusion equation, a popular framework based on fractional calculus [48]. The fractional diffusion equation was previously investigated in the context of random-walk theory, and it describes a Lévy flight process [48]. Here, we have provided a microscopic justification for it. The fractional diffusion equation [48,71,72]

\[ \frac{\partial}{\partial t} P(x, t) = K_x \nabla^\beta P(x, t) \]  
(105)

was the phenomenological starting point for the description of the experiments in the work of Sagi et al. [3]. Our work (see also Ref. [26]) provides the exponent \( \nu \), Eqs. (7) and (30), in terms of the recoil energy and lattice depth, \( \beta = 1 \), and the anomalous diffusion coefficient \( K_x \) from Eq. (76). Indeed, the experiment [3] found the value \( \beta = 1 \), so the time derivative on the left-hand side is a first-order time
derivative. The fractional space derivative $\nabla^p$ is a Weyl-Reitz fractional derivative [48]. To see the connection between our results and the fractional equation, we reexpress Eq. (74) as

$$u\tilde{P}(k, u) - 1 = -K_x |k|^p \tilde{P}(k, u),$$

(106)

which is the Fourier-Laplace transform of Eq. (105). Here, we recall [48] that the Fourier-space representation of $\nabla^p$ is $-|k|^p$ and that $u\tilde{P}(k, u) - 1$ is the representation of the time derivative in Laplace-Fourier space of a $\delta$-function initial condition centered on the origin. We see that $P(x, t)$ for $2/3 < \nu < 2$ satisfies the fractional diffusion equation. In other words, for initial conditions starting on the origin, the solution of Eq. (105) is the Lévy PDF, Eq. (78).

However, the use of the fractional diffusion equation must be performed with care. It predicts a seemingly unphysical behavior: The mean-square displacement is apparently infinite. Although, as mentioned in the Introduction, the mean-square displacement was shown to exhibit superdiffusion, experimentally by Katori et al. [1] and in simulations in [2], it is nevertheless finite at all finite times. In fact, the coupling between excursion duration $\tau$ and extent $\chi$ analyzed herein implies that the fractional equation is valid only in a scaling region, $|x| < r^{1/2}$, and thus the Lévy distributions obtained analytically here and used phenomenologically in the Weizmann experiment [3] describe the center part of the spreading distribution of the particles, but the tails exhibit a cutoff for $x > r^{1/2}$. In other words, the Lévy distribution describing the center part of the packet for $1/5 < D < 1$ does not contain information on the correlations, and to experimentally investigate correlations in this regime, one must probe the tails of the packet.

The renewal approach is the basis of the coupled CTRW theory developed here. Renewal theory turned out to be predictive in the sense that, while the set of zero crossings is nontrivial in the continuum limit, we could avoid this mathematical obstacle by introducing modified paths with a $\pm \epsilon$ starting point after each zero hitting. Physical observables, like $\tau$, $P(x, t)$, $K_x$, and $(\chi^2)$, do not depend on $\epsilon$ in the limit, as expected. Technically, this is due to a cancellation of $\epsilon$. For example, both $\langle \tau \rangle$ [Eq. (57)] and $q^\star$ [Eq. (66)] depend linearly on $\epsilon$; hence, their ratio is $\epsilon$ independent, and the anomalous diffusion coefficient $K_x$ [Eq. (76)] becomes insensitive to $\epsilon$. The same is true for the normal diffusion constant $K_2$ [Eq. (102)] and other physical observables [Eqs. (69), (89), (92), and (93)]. A close reading of Appendixes A and B will reveal that this cancellation is not exactly trivial. In our work, we use the \textit{regularized} Bessel process for the calculation of the PDFs of the first passage time $g_t(\tau)$ and jump lengths $q(\chi)$. Hence, for the calculations, the details of the potential $V(p)$ for small $p$ are important [but not for $\mathcal{B}(\nu, 2)$]. If we use instead the PDFs $g_{\nu}(\tau)$ and $q_{\nu}(\chi)$ (see Appendixes A and B) of the nonregularized process, i.e., a logarithmic potential for any $p$, some (but not all) of the $\epsilon$ cancellation will not take place [see Eq. (67) where both $q^\star \propto \epsilon$ and $g^\star \propto \epsilon$ for the regularized process, which allows for the cancellation]. For example, both for the regularized and the nonregularized processes, $g_t(\tau) \sim g^\star \tau^{-(1+\alpha)}$ (so the exponent $\alpha$ is identical in both cases); however, $g^\star \propto \epsilon$ for the regularized process, while $g^\star \propto \epsilon^{2\nu}$ for the nonregularized case, Eq. (A7). The cancellation of $\epsilon$ is further discussed in Appendix F for three types of random durations. That appendix shows that, while the first passage time PDF $g_t(\tau)$ depends on $\epsilon$ (see Appendix A), the time interval straddling time $t$ and the backward and forward recurrence times are insensitive to this parameter. Thus, one upshot of Appendix F is a further demonstration that this intuitive $\epsilon$ renewal trick actually works. A rigorous mathematical treatment of the problem will lead to a stronger foundation of our results.

Regularization of the process is crucially important in the Lévy and Gaussian phase $D < 1$ where observables depend on the details of the potential $V(p)$. This is related to work of Martin et al. [41] on the classification of boundaries for the nonregularized Bessel process. They showed that, using our notation, for $D < 1$, $p = 0$ is an exit boundary while for $D > 1$, it is a regular boundary ($D < 0$ corresponds to an entrance boundary condition, which is not relevant to our work). For an exit boundary starting on $p = 0$, it is impossible to reach a finite-momentum state $p$, so clearly we cannot allow such a boundary in our physical problem. For that reason, we need to consider the regularized first passage time problem, at least when $D < 1$, and then the boundary is regular. A regular boundary on $p = 0$ means the diffusion process can enter and leave from the boundary. Therefore, for $D > 1$, our final results do not depend on the shape of $V(p)$, besides its asymptotic limit; namely, one may replace the regularized Bessel process with the nonregularized one. To see this, notice that $\langle x^2 \rangle$ for $D > 1$ depends on $\alpha < 1$ but not on small-scale properties of $V(p)$, and since $\alpha$ is the first passage exponent for both the regularized and nonregularized processes, it does not really matter which process is the starting point of the calculation. Further evidence for such a behavior is in Appendix F, where statistics of durations for $D > 1$ are shown to depend on $\alpha < 1$ but not on the shape of $V(p)$. Thus, $D = 1$ marks the transition from Richardson to Lévy behavior, the diverging of the partition function $Z$, namely, the equilibrium velocity distribution, turns nonnormalizable, and according to Martin et al., the boundary of the Bessel process switches from regular to exit.

Additional theoretical work is needed on the leakage of particles and, more generally, evaporation (see more details below). Assuming that energetic particles are those which get evaporated, and assuming that this takes place through the boundary of the system, one would be interested in the first passage time properties of particles, from the center of
the system to one of its boundaries. Previous theoretical work on first passage times for Lévy walks and flights might give a useful first insight [73–77]; however, it is left for future work to compare these simplified models with the microscopic semiclassical picture of the underlying dynamics. Evaporation is an important ingredient of cooling (beyond the Sisyphus cooling) since it gets rid of the very energetic particles; hence, this line of research has practical applications. Theoretically, understanding the boundary conditions for fractional diffusion equations, needed for the calculation of statistics of first passage times, is still an open challenge. Yet another challenge is to characterize the joint PDF $W(x, p, t)$ and, in particular, the correlations between position and momentum, which are expected to be nontrivial. Again, some elegant ideas and tools were recently developed [78] within the Lévy walk framework; however, more work is needed for direct comparison with the physical picture under consideration.

B. Cold-atom experiments

Experiments on ultracold diffusing particles in optical lattices have been performed on a single $^{24}$Mg$^+$ ion [1] and an ensemble of Rb atoms [2,3,20]. Experiments on the spreading of the density of atoms yield ensemble averages like the mean-square displacement and the density $P(x, t)$. Analysis of individual trajectories yields, in principle, deep insight on paths, for example, information on the first passage time in velocity space, i.e., $g(\tau)$, zero crossing times, statistics of excursions, meanders, etc. For a comprehensive understanding of the dynamics of the particles, both types of experiments are needed. The theory presented herein provides statistical information on the stochastic trajectories, e.g., zero crossing events, which in turn are related to the spreading of the packet of particles. Thus, from single trajectories, one may estimate $g(\tau)$, $B(v_{1/2})$, and $\psi_j(\chi, \tau)$. From $g(\tau)$, one may then obtain the exponent $\alpha$, which can be used to predict the qualitative features of the spreading of the ensemble of particles. For example, with an estimate of $\alpha$, we can determine the phase of motion, be it Richardson, Lévy, or Gauss. Of course, a more quantitative investigation is now possible since we have analytically related $\alpha$, the mean-square displacement, and $K_\nu$ with microscopical parameters like the optical potential depth $U_0$. In principle, statistics of Bessel excursions, so far of great interest mainly in the mathematical literature, could be detected in single-particle experiments. These single-ion experiments are also ideal in the sense that collisions/interparticle interactions do not play a role, and they also provide insights on ergodicity. Three specific unsolved issues are as follows:

(i) Can the Richardson phase be experimentally demonstrated? So far, this phase was obtained in our theory and Monte Carlo simulations [2], while experiments exhibit ballistic diffusion as the upper limit [1–3]. This might be related to the sublattice limit of taking time to infinity before the depth of the optical lattice approaches zero, since clearly in the absence of an optical lattice, the fastest particles are ballistic.

(ii) While our theory correctly predicts a Lévy phase in agreement with experiments, in Ref. [3] two exponents were used to fit the data. In contrast, the theory we developed suggests a single Lévy exponent $\nu$. This could be due to the leakage of particles and also to the $r^{3/2}$ cutoff of the Lévy density, which are due to the correlations between $\chi$ and $\tau$ investigated in this manuscript. Avoiding the leakage of particles is an experimental challenge, and once this is accomplished, a more informed comparison of our theoretical prediction on the Lévy phase with a single characteristic exponent $\nu$ could be made. Success on this front would provide an elegant demonstration of Lévy’s central-limit theorem, with the characteristic exponent controlled by the depth of the optical potential.

(iii) As pointed out in Ref. [2], losses of atoms lead to an underestimate of the diffusion exponent, as many superdiffusing atoms are lost. Characterization of these losses both theoretically and experimentally could advance the field, since this yields insight on the underlying processes and leads us towards better control of the particles. Specifically, we do not know how many particles are kept in the system, versus time, for varying strength of the optical potential.

(iv) In the Lévy phase, measurements of noninteger moments $\langle |x|^q \rangle$ with $q < \nu < 2$ will, according to theory, exhibit Lévy scaling $\langle |x|^q \rangle \sim r^{\nu q}$. In contrast, the main focus of experiments so far was the second moment $q = 2$, which is difficult to determine statistically since, as we have shown here, it depends on the tails of $P(x, t)$ and very fast particles. Thus, ideally, a wide spectrum of moments $\langle |x|^q \rangle$ should be recorded in experiment, with the low-order moments $q < \nu$ giving information on the center part of the density and the higher-order moments giving information on the correlations and tails.

In the case where experiment and theory do not reconcile, we will have a strong indication that the current semiclassical theory is not sufficient; then, we will be forced to investigate at least four other aspects of the problem:

(a) Effects of collisions on anomalous diffusion of the atoms.

(b) Effect of higher dimensions.

(c) Quantum effects beyond the semiclassical approach used here. In particular, it would be very interesting to simulate this system with full quantum Monte Carlo simulations [6], to compare the semiclassical theory with quantum dynamics. We note that the Richardson phase, which, as mentioned, was not observed in experiments, is actually a heating phase, and quantum simulations become difficult because the numerical lattice introduces a cutoff on velocities which induces artificial ballistic motion.

(d) Other cutoff effects that modify the anomalous diffusion. For example, at high enough velocities, Doppler cooling is expected to diminish the fast particles.
In Ref. [79], we simulated the effect of Doppler friction and showed that the anomalous character of the diffusion is kept unchanged, at least for a certain reasonable set of parameters. However, a general rule on the influence of Doppler cooling has not yet been established, and an experimentalist with a specific set of parameters in mind might wish to numerically test the magnitude of this effect on the anomalous spreading.

An interesting approach was recently suggested by Dechant and Lutz [22]. They investigated the multifractal nature of the moments $\langle |x|^q \rangle$ of the process and considered initial conditions different than ours. They assumed that in the first stage, of duration $t_c$, the particles are cooled in a confining field (which inhibits the spreading). Then, the momentum of particles relaxes to a state described by the infinite covariant density [25], which depends on $t_c$. The particles are then released, and their spreading is recorded for a duration $t$. We consider the case $t \gg t_c$, while Ref. [22] considered the opposite case. The experiment in Ref. [3] is conducted under the conditions investigated here, namely, that the spreading time is much longer than the preparation time. As shown by Hirschberg et al., starting with power-law distributions will dramatically influence the spreading, both in momentum space [24] and in space. Indeed, large $t_c$ implies power-law initial conditions, Eq. (9), with a $t_c$-dependent cutoff [25]. In that sense, diffusivity is sensitive to the initial preparation of the system, and an experimental verification of these effects would indicate the fundamental difference between transport in these systems and normal transport, which does not depend on initial conditions. This is clearly related to the strong sensitivity we have found of the mean-square displacement on a single jump event, described by the Bessel meander.

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APPENDIX A: THE WAITING-TIME PDF

The waiting time is the time it takes the particle starting with momentum $p_i > 0$ to reach $p_f < p_i$ for the first time. Here, we investigate its PDF, $g(t)$.\[\text{\textbf{A1}}\]

namely, the probability that a particle initially at $p_i$ does not cross the boundary $p_f < p_i$ in the time interval whose length is $\tau$, satisfies\[\text{\textbf{A1}}\]

Here, $S = 1$ for $\tau = 0$ since the particle’s escape time cannot be zero. Furthermore, $S \to 0$ when $p_i \to p_f$, since in that case, the particle starts out at the boundary, and $S = 1$ if one starts at $p_i = \infty$. The random time it takes a particle starting on $p_i$ to reach $p_f$ for the first time is $\tau$, and its PDF is $g_\tau(\tau) = -\partial_\tau S(\tau)$. Here, the subscript $B$ denotes the Bessel process. Since $S(\tau)|_{\tau=0} = 1$, we have, in Laplace $\tau \to u$ space, the following simple relation:

Using Eq. (A2) and the Laplace transform of Eq. (A1), we find

Note that $\hat{g}_\tau(u)|_{p_i=p_f} = 1$ since $g_\tau(\tau)|_{p_i=p_f} = \delta(\tau)$. The solution of Eq. (A3) is [80]

Here, $K_\alpha(\cdot)$ is the modified Bessel function of the second kind, and as before, $\alpha = 1/2 + 1/(2D)$.

If $\alpha > 1$, the small $u$ expansion of Eq. (A4) yields $\hat{g}_\tau(u) \sim 1 - u\langle \tau_B \rangle + \cdots$, where $\langle \tau_B \rangle$ is the average first passage time. Expanding Eq. (A4), we find

Notice that $\langle \tau_B \rangle$ diverges when $\alpha \to 1$, corresponding to $D \to 1$. Not surprisingly, the average time in Eq. (A5) is generally different from the expression for the average waiting time for the regularized process, Eqs. (56) and (57).

Expanding Eq. (A4) for small $u$ for the case $0 < \alpha < 1$, we get

Inverting to the time domain, we find, using a Tauberian theorem, the long-time behavior of the PDF,

\begin{align}
\hat{g}_\tau(u) &= \frac{1}{\sqrt{2\pi}} \frac{\Gamma(1-\alpha)}{\Gamma(1+\alpha)} \left(\frac{1}{2}\right)^{2\alpha} ((p_i)^{2\alpha} - (p_f)^{2\alpha}) \frac{u^{\alpha}}{D} \cdots. \\
\text{\textbf{A6}}
\end{align}

We first briefly investigate the first passage time problem for the Bessel process following Ref. [35]. The Bessel process corresponds to the case $F(p) = -1/p$, so the force diverges at the origin. In the next subsection, we will consider the regularized force $F(p) = -p/(1 + p^2)$. According to Ref. [50], the survival probability $S(\tau)$,
One can show that this fat-tail behavior is also valid for the regime \( \alpha > 1 \). The procedure involves expansion of Eq. (A4) beyond the first two leading terms; for example, for \( 1 < \alpha < 2 \), one finds \( \hat{g}_p(u) = 1 - u(R_0) + Cu^\alpha \), and the third term gives the tail of the PDF. Note that the PDF \( g_p(\tau) \) yields the same exponent as in Eq. (11). However, to calculate the amplitude \( g' \), we must consider the regularized process.

### 2. First passage time for the regularized process

For the optical lattice problem, we need to treat the regularized force \( F(p) = -p/(1 + p^2) \). Our aim is to find the asymptotic behavior in Eq. (11), while the average waiting time is given already in Eqs. (56) and (57). From these, we know that for \( \alpha < 1 \) the average first passage time from \( p_i \) to \( p_f \) is infinite (the derivation can be easily generalized for arbitrary initial and final states). Furthermore, for long times, the large-momentum behavior of \( F(p) \) plays the crucial role, and hence we expect, for small \( u \),

\[
\hat{g}(u) \approx 1 - Gu^\alpha; \quad 0 < \alpha < 1, \quad (A8)
\]

where we must determine \( G \), which then gives the amplitude \( g' \) introduced in Eq. (11). The equation for \( \hat{g}(u) \) is the same as Eq. (A3) but with the regularized force

\[
D(\partial_p)^2 \hat{g}(u) - \frac{p_f}{1 + (p_f)^2} \partial_p \hat{g}(u) - u \hat{g}(u) = 0. \quad (A9)
\]

We need to solve for \( \hat{g}(u) \) for small \( u \), so to leading order, we drop the last term and find

\[
D(\partial_p)^2 \hat{g}(u) + F(p_i)\partial_p \hat{g}(u) = 0. \quad (A10)
\]

Hence,

\[
\hat{g}(u) = C_1(u) \int_{p_f}^{p_i} e^{V(p)/D} dp' + C_2(u), \quad (A11)
\]

where \( V(p) = \ln(1 + p^2)^{1/2} \). Since \( \hat{g}(u) = 1 \) when \( p_i = p_f = 1 \), \( C_2(u) = 1 \). This approximation breaks down, however, when \( p_i \) is too large since then \( \hat{g}(u) \propto p_i^{-1/2} \), and the last term in Eq. (A9) is comparable in size to the first two when \( p_i \sim u^{-1/2} \). In the large \( p_i \) regime, however, \( p_f \gg 1 \), so we can approximate \( F(p_i) \) by its Bessel form, \( F(p_i) \approx -1/p_i \), and the solution is, as above,

\[
\hat{g}(u) \approx C_3(u)K_\alpha \left( p_f \sqrt{\frac{u}{D}} \right)^{p_f}. \quad (A12)
\]

where we cannot use the \( p_f \to p_i \) limit to normalize \( \hat{g}(u) \) since \( p_f \) is not necessarily large. Nevertheless, \( \hat{g}(u = 0) = 1 \) and \( K_\alpha(z) \sim (\Gamma(\alpha)/2)(z/2)^{-\alpha} \) for \( z \to 0 \), implying that to leading order in \( u \). For \( 1 \ll p_i \ll u^{-1/2} \), our two approximations must agree, and so we find, using the second-order expansion of \( K_\alpha(z) \),

\[
1 - \frac{\Gamma(1 - \alpha)}{\Gamma(1 + \alpha)} \left( p_f \sqrt{\frac{u}{D}} \right)^{2\alpha} \approx 1 + C_1(u) \frac{p_i^{1/2 + 1}}{1/D + 1}. \quad (A14)
\]

and recognizing that \( (1/D + 1) = 2\alpha \),

\[
C_1(u) \approx \frac{2\alpha}{\Gamma(1 + \alpha)} \left( \frac{u}{4D} \right)^{\alpha}. \quad (A15)
\]

Thus, using Eq. (A11),

\[
\hat{g}(u) \approx 1 - \left[ \frac{2\alpha}{\Gamma(1 + \alpha)} \left( \frac{1}{4D} \right)^{\alpha} \int_{p_f}^{p_i} e^{V(p)/D} dp' \right] u^\alpha. \quad (A16)
\]

Using the Tauberian theorem, which implies \( u^\alpha \to r^{-(1 + \alpha)/\Gamma(-\alpha)} \), we find the large \( r \) behavior

\[
g(r) \sim \frac{2\alpha}{\Gamma(1 + \alpha)} \left( \frac{1}{4D} \right)^{\alpha} r^{-(1 + \alpha)}. \quad (A17)
\]

It can be shown, with additional work, that this result is also valid for \( \alpha > 1 \). To conclude, we find \( p_i = \epsilon \ll 1 \) for the initial state,

\[
g^* = \frac{2\alpha}{(4D)^{\alpha} \Gamma(1 + \alpha)} \int_0^\epsilon e^{V(p)/D} dp' \approx \frac{2\alpha}{(4D)^{\alpha} \Gamma(1 + \alpha)}. \quad (A18)
\]

This, of course, vanishes as \( \epsilon \to 0 \), in this case linearly, as opposed to the \( p_i \to \epsilon, p_f \to 0 \) limit of Eq. (A7), where the dependence is of higher order. As shown in this paper, for the purpose of calculation of the asymptotic behavior of \( P(x, t) \), all we need is \( \partial_\tau g^* \) when \( \epsilon \to 0 \), which is a finite constant that depends on \( D \) only. The full shape of \( V(p) \) is unimportant for the calculation of \( g^* \), indicating a degree of universality. In contrast, \( \langle \tau \rangle \) depends on \( Z \) and so is nonuniversal. Since \( g^* \) is a measure of the long-time behavior, the detailed shape of the potential is not important, provided it is regularized.

### 3. Moments of \( g(\tau) \)

In the text, we show that \( \langle \tau \rangle \) and \( \langle \tau^2 \rangle \) are relevant when \( D < 1/5 \), so these averages do not diverge. The moments \( \langle \tau \rangle \) and \( \langle \tau^2 \rangle \) are found using

\[
\hat{g}(u) = 1 - u(\tau) + u^2(\tau^2)/2 - u^3(\tau^3)/6 \ldots. \quad (A19)
\]
we get
\[
\left[ D(\partial_p)^2 - \frac{p}{1 + p^2} \partial_p - u \right] \hat{g}(u) = 0, \tag{A20}
\]
we get
\[
\left[ D(\partial_p)^2 - \frac{\tau - p}{1 + \tau p^2} \partial_p \right] (\tau) = -1. \tag{A21}
\]
\[
\left[ D(\partial_p)^2 - \frac{p}{1 + p^2} \partial_p \right] (\tau^2) = -2(\tau). \tag{A22}
\]
\[
\left[ D(\partial_p)^2 - \frac{p}{1 + p^2} \partial_p \right] (\tau^3) = -3(\tau^2). \tag{A23}
\]
Solving Eq. (A21),
\[
\langle \tau \rangle = \frac{1}{D} \int_0^p dy v(y)/D \int_{\nu}^\infty e^{-v(x)/D} dx, \tag{A24}
\]
so \( \langle \tau \rangle = 0 \) when \( p = 0 \), as expected for the first passage from \( p \) to the origin. For the second and third moments,
\[
\langle \tau^2 \rangle = \frac{2}{D} \int_0^p dy v(y)/D \int_{\nu}^\infty e^{-v(x)/D} \langle \tau(x) \rangle dx \tag{A25}
\]
and
\[
\langle \tau^3 \rangle = \frac{3}{D} \int_0^p dy v(y)/D \int_{\nu}^\infty e^{-v(x)/D} \langle \tau^2(x) \rangle dx. \tag{A26}
\]
Thus, in the \( p = \epsilon \to 0 \) limit,
\[
\langle \tau^3 \rangle \sim \frac{3\epsilon}{D} \int_{0}^{\infty} e^{-v(x)/D} \langle \tau^2(z) \rangle dz. \tag{A27}
\]

\section*{APPENDIX B: THE JUMP LENGTH PDF \( q(x) \)}

The excursion length \( x \) is the distance the particle travels from its beginning with momentum \( p_i > 0 \) until it reaches the momentum origin \( p_f = 0 \) for the first time. Here, we investigate its PDF \( q(x) \), which in the limit of \( p_i \to 0 \) and for the regularized force field, gives the jump length PDF. Previously, this PDF was investigated for the \(-1/p\) force field in Ref. [6]. Note that, in what follows, \( p_i > 0 \) so \( x > 0 \). As elsewhere in this paper, from symmetry, positive and negative \( x \) are equally probable, so we may restrict our attention to \( x > 0 \).

1. PDF of jump lengths for the Bessel process

As mentioned, for the Bessel process, \( F(p) = -1/p \). The PDF \( q_B(x) \), which depends of course on the start and end points \( p_i \) and \( p_f \), satisfies the following backward equation [50]:
\[
D(\partial_p)^2 q_B - \left( 1/p_i \right) \partial_p q_B - p_i \partial_x q_B = 0, \tag{B1}
\]
where the subscript \( B \) again stands for Bessel. Since \( x > 0 \), we define the Laplace transform
\[
\hat{q}_B(s) = \int_0^\infty \exp(-sx) q_B(x) dx. \tag{B2}
\]
When \( p_f \to p_i \), we have \( q_B(x) \to \delta(x) \) and \( q_B(x)_{x=0} = 0 \) if \( p_i \neq p_f \) since it takes time for the particle to reach the boundary, and hence we cannot get an excursion whose size is zero. Of course, in the opposite limit of large \( x \), \( \lim_{x \to \infty} q_B(x) = 0 \). In Laplace space, Eq. (B1) yields
\[
(\partial_p)^2 \hat{q}_B(s) - (1/Dp_i) \partial_p \hat{q}_B(s) - (p_i s/D) \hat{q}_B(s) = 0. \tag{B3}
\]
It is easy to verify that the appropriate solution is
\[
\hat{q}_B(s) = \frac{\left(p_i \right)^{3\nu/2} K_\nu \left( \sqrt{D} p_i \right)^{3/2}}{\left(p_i \right)^{\nu/2} K_\nu \left( \sqrt{D} p_f \right)^{3/2}}, \tag{B4}
\]
with \( \nu = (1 + D)/3D \), and as in the previous appendix, \( K_\nu \) is the modified Bessel function of the second kind.

Our goal is to find the large \( x \) behavior of \( q_B(x) \), so we expand Eq. (B4) in the small \( s \) limit, finding, for \( 0 < \nu < 1 \),
\[
\hat{q}_B(s) \sim 1 - \frac{1}{\nu} \left[ (p_i)^{3\nu} - (p_f)^{3\nu} \right] \frac{\Gamma(1-\nu)}{\Gamma(1+\nu)} (3\nu-1)s^{3\nu-1} \cdots. \tag{B5}
\]
The first term on the left-hand side is simply the normalization, and the \( s^\nu \) term indicates that the average excursion length diverges with \( s \) for large \( x \),
\[
q_B(x) \sim \left(\frac{1}{3}\right)^{2\nu} \frac{(p_i)^{3\nu} - (p_f)^{3\nu}}{\Gamma(\nu)} (3\nu-1)x^{3\nu-1}. \tag{B6}
\]
With a similar method, one can show that Eq. (B6) also holds for \( 1 < \nu < 2 \). There, the expansion, Eq. (B5), contains three terms; the additional term yields the average of \( x \), which is now finite. In Ref. [6], a more complicated method was used to investigate the same problem, and a similar result was found, although with a typographical error. That reference reports \( (3\nu+1)^s \) in the prefactor, whereas we find \((3\nu-1)^s\), a difference with some importance since the prefactor found here goes to zero when...
$D \to \infty$, which is necessary to obtain reasonable physical results.

2. PDF of $\chi$ for the regularized process

To obtain the large $\nu$ behavior of $q(\chi)$ for the regularized force $F(p) = -p(1 + p^2)$, we follow the same steps carried out in Appendix A 2. The equation to solve is

$$D(\partial_{p_i})^2 \hat{q} + F(p_i) \partial_{p_i} \hat{q} - p_i \partial_{q} \hat{q} = 0.$$ \hspace{1cm} (B7)

For $\nu < 1$, we switch to Laplace space $\chi \to s$, i.e., $D(\partial_{p_i})^2 \hat{q} + F(p_i) \partial_{p_i} \hat{q} - p_i \hat{q} = 0$, and drop the last term for small $s$, yielding

$$\hat{q}(s) \approx C_i(s) \int_{p_i}^{p_i+\epsilon} e^{\nu(p')/D} dp' \hspace{1cm} (B8)$$

after applying the boundary condition $\hat{q}\big|_{p_i+\epsilon} = 1$. Again, for $p_i \gg 1$, we can approximate $F(p_i)$ by its large $p_i$ approximation, yielding

$$\hat{q}(s) \approx C_3(s) p_i^{3\nu/2} K_s \left( \frac{2sp_i^2}{3D} \right). \hspace{1cm} (B9)$$

In the small $s$ limit, $q(s \to 0) = 1$ implies

$$C_3(s) \approx \frac{2}{\Gamma(\nu)} \left( \frac{1}{3} \sqrt{\frac{s}{D}} \right)^\nu \hspace{1cm} (B10)$$

to leading order. These two approximations must agree for $1 \ll p_i \ll s^{-1/3}$, yielding

$$C_4(s) \approx 3^{1-2\nu}(3\nu - 1) \frac{\Gamma(1 - \nu)}{\Gamma(\nu)} \nu \frac{1}{\sqrt{s}}. \hspace{1cm} (B11)$$

This implies that

$$q(\chi) \approx \frac{1}{3^{2\nu-1} \Gamma(\nu)} \nu \int_{p_i}^{p_i+\epsilon} e^{\nu(p')/D} dp' \chi^{-1-\nu}. \hspace{1cm} (B12)$$

This calculation was done assuming $p_i, p_f$ are positive. Allowing also for negative momenta, and so negative $\chi$, we get an additional prefactor of $1/2$:

$$q(\chi) \approx \frac{1}{3^{2\nu-1} \Gamma(\nu)} \nu \int_{p_i}^{p_i+\epsilon} e^{\nu(p')/D} dp' |\chi|^{-1-\nu}. \hspace{1cm} (B13)$$

One can show that this result is valid in the whole domain of interest $1/3 < \nu < 2$, i.e., the domain where the variance of $\chi$ is infinite.

3. The variance of $\chi$

Here, we obtain the finite $\langle \chi^2 \rangle$ for the case $\nu > 2$ for the regularized force $F(p) = -p(1 + p^2)$. Using the Laplace $\chi \to s$ transform of Eq. (B7), we find the following backward equation for $\hat{q}(s)$:

$$D \frac{\partial^2}{\partial (p_i)^2} \hat{q}(s) + F(p_i) \frac{\partial}{\partial p_i} \hat{q}(s) - p_i s \hat{q}(s) = 0. \hspace{1cm} (B14)$$

Here, we used $q(\chi)|_{p=0} = 0$ since the particle starting with $p_i > 0$ cannot reach zero momentum without traveling some finite distance. The Laplace transform $\hat{q}(s)$ is expanded in $s$:

$$\hat{q}(s) = \int_0^\infty e^{-st} q(\chi) d\chi = 1 - s \langle \chi \rangle + s^2 \langle \chi^2 \rangle/2 + \cdots. \hspace{1cm} (B15)$$

Here, $\langle \chi \rangle$ is the average jump size, for positive excursions, i.e., those that start with $p_i > 0$, and $\langle \chi^2 \rangle$ is the second moment. Note that in the original model, we have both positive and negative excursions, and so we have $p_i = +\epsilon$ or $p_i = -\epsilon$ with probability $1/2$; hence, for that case $-\infty < \chi < \infty$ and from symmetry $\langle \chi \rangle = 0$. The variance of the original process is $\langle \chi^2 \rangle$ because of symmetry. As we restrict ourselves to $p_i > 0$, $\chi > 0$ and $\langle \chi \rangle$ in Eq. (B15) is finite.

Inserting Eq. (B15) into Eq. (B14), we find

$$D \frac{\partial^2}{\partial (p_i)^2} \left[-s \langle \chi \rangle + s^2 \langle \chi^2 \rangle/2 + \cdots \right] - \frac{p_i}{1 + (p_i)^2} \partial_{p_i} \left[-s \langle \chi \rangle + s^2 \langle \chi^2 \rangle/2 + \cdots \right] - p_i s \left[1 - s \langle \chi \rangle + s^2 \langle \chi^2 \rangle/2 + \cdots \right] = 0. \hspace{1cm} (B16)$$

The $s^1$ terms give

$$-D \frac{\partial^2}{\partial (p_i)^2} \langle \chi \rangle + \frac{p_i}{1 + (p_i)^2} \partial_{p_i} \langle \chi \rangle - p_i = 0, \hspace{1cm} (B17)$$

while the $s^2$ terms are

$$2D \frac{\partial^2}{\partial (p_i)^2} \langle \chi^2 \rangle - \frac{1}{2} \frac{p_i}{1 + (p_i)^2} \partial_{p_i} \langle \chi^2 \rangle + p_i \langle \chi \rangle = 0. \hspace{1cm} (B18)$$

This means that the first and then the second moment can be found by repeated integration. The boundary conditions are that both the first and second moments are zero if $p_i = 0$, while they diverge if $p_i = \infty$. Using a reflecting boundary at $p_i = \infty$, we find (see, e.g., Ref. [50], Chapter 5)
where the lower bound is the momentum origin and the potential is given in Eq. (8). Integrating Eq. (B18), we find

$$\langle \chi^2 \rangle = \frac{2}{D} \int_0^\infty \text{d}z e^{-\frac{V(z)}{D}} \int_0^\infty \text{d}z e^{-\frac{V(z)}{D}} \langle \chi(z) \rangle. \quad (B20)$$

As explained in the text, we consider only the limit when \( p_i = \epsilon \) is small, so we have

$$\langle \chi^2 \rangle \sim \frac{2\epsilon}{D} \int_0^\infty \text{d}z e^{-\frac{V(z)}{D}} \int_0^\infty \text{d}z e^{-\frac{V(z)}{D}} \langle \chi(z) \rangle. \quad (B21)$$

since \( V(0) = 0 \). Using Eq. (B19) and integrating by parts, we get

$$\langle \chi^2 \rangle \sim \frac{2\epsilon}{D^2} \int_0^\infty \text{d}z e^{-\frac{V(z)}{D}} \left[ \int_0^\infty \text{d}z e^{-\frac{V(z)}{D}} \right]^2. \quad (B22)$$

Note that from symmetry, \( V(p) = V(-p) \), so we have

$$\int_0^\infty \text{d}z \exp\left[-\frac{V(z)}{D}\right] = -\int_\infty^0 \text{d}z \exp\left[-\frac{V(z)}{D}\right].$$

**APPENDIX C: EXPANSION OF EQ. (62)**

Since \( g(\tau) \approx g^* \tau^{-3/2-\gamma} \) for large \( \tau \), with \( \tau = 1/(2D) \), we assume that there exists some \( t_0 \) such that \( g(\tau) = g^* \tau^{-3/2+\gamma} \) for \( \tau > t_0 \). Then, using Eq. (59) and \( \nu = 1/3 + 2\beta/3 \),

$$\psi(k, u)_{\kappa = 0} = 1 - \int_0^t \text{d}\tau \left[ 1 - B(k \sqrt{D \tau^2}) \right] g(\tau)$$

$$+ \int_t^\infty \text{d}\tau \left[ 1 - B(k \sqrt{D \tau^2}) \right] g(\tau). \quad (C1)$$

The first term on the right-hand side yields the normalizing condition \( \psi(k, u)_{\kappa = 0} = 1 \); the second term is zero for \( k = 0 \), and because of the symmetry of the jumps, \( B(v_{3/2}) = B(-v_{3/2}) \), this term, when expanded in \( k \), gives a \( k^2 \) term. As we will show, the last term gives a \( |k|^\nu \) term, so as long as \( \nu < 2 \), we can neglect the second term. Then, changing variables to \( k = k \sqrt{D \tau^2} \), we find

$$\psi(k, u)_{\kappa = 0} = 1 - \frac{2}{\beta} g_\nu (Dk^2)^{-\nu/2} \int_0^\infty \frac{1 - B(k \sqrt{D \tau^2})}{(k \sqrt{D \tau^2})^{1+\nu}} \text{d}\tau. \quad (C2)$$

Here, we have taken the small \( k \) limit in such a way that \( k(t_0)^{3/2} \), which appears in the lower limit of an integral, is negligible.

From Eq. (C2), we must investigate the integral

$$I_\nu = \int_0^\infty \frac{1 - B(k \sqrt{D \tau^2})}{(k \sqrt{D \tau^2})^{1+\nu}} \text{d}\tau. \quad (C3)$$

From symmetry, \( B(v_{3/2}) = B(-v_{3/2}) \), and hence the Fourier transform \( \tilde{B}(k) \) is a real function. Thus,

$$I_\nu = \text{Re} \left[ \int_0^\infty \frac{1 - \tilde{B}(k \sqrt{D \tau^2})}{(k \sqrt{D \tau^2})^{1+\nu}} \text{d}\tau \right]. \quad (C4)$$

and \( k_\nu \) will eventually be taken to zero. Using the definition of the Fourier transform,

$$I_\nu = \kappa \tilde{B}(k \sqrt{D \tau^2} - 2 \int_0^\infty \text{d}v \tilde{B}(v_{3/2}) \int_k^\infty \frac{e^{i(k \sqrt{D \tau^2}/k^{1+\nu})}}{k^{1+\nu}} \text{d}k]. \quad (C5)$$

Let \( iv_{3/2} = -x \), and then using integration by parts,

$$\text{Re} \left[ \int_k^\infty \frac{e^{i(k \sqrt{D \tau^2}/k^{1+\nu})}}{k^{1+\nu}} \text{d}k \right] = \text{Re} \left[ e^{-k \sqrt{D \tau^2} / k^{1+\nu}} - x \sqrt{1 - 1/\nu} \right]. \quad (C6)$$

Inserting the last expression in Eq. (C5), the diverging \( (k_\nu)^{-\nu} \) terms cancel each other, and we find

$$I_\nu = \frac{2\Gamma(1-\nu)}{\nu} \text{Re} \left[ \int_0^\infty e^{i(k \sqrt{D \tau^2}/k^{1+\nu})} \text{d}k \right]. \quad (C7)$$

Using \( x = -iv_{3/2} \), we get

$$I_\nu = -\langle |v_{3/2}|^\nu \rangle \sin \left( \frac{\pi(1+\nu)}{2} \right) \Gamma(-\nu), \quad (C8)$$

where \( \langle |v_{3/2}|^\nu \rangle = \int_0^\infty |v_{3/2}|^\nu B(v_{3/2}) \text{d}v_{3/2} \). Inserting Eq. (C8) into (C2) and using (C3), we get Eq. (63).

**APPENDIX D: ON THE SIMULATIONS**

1. The simulation

In this appendix, we briefly discuss the simulations presented in the paper to test our predictions. We have two classes of simulations: One generates Bessel excursions, and the other solves the Langevin equations. The first is based on a discrete random-walk treatment of the excursion, wherein the particle takes a biased walk to the right or left in momentum space every time step. The degree of bias varies with \( p \) in accordance with the force \( F(p) \). Our dimensionless control parameter \( D \) is the ratio of the coefficient of the diffusion term to the coefficient of the \( 1/p \) large-\( p \) behavior of the force. Since for our simple random walk the diffusion constant is \( 1/2 \), the parameter \( D \) enters via the strength of the bias. The probability to move right or left is then given by

$$P_{\pm}(p) = \frac{1}{2} \left( 1 \pm \frac{1}{2D} F(p) \Delta p \right). \quad (D1)$$
of unity, it is sufficient to take \( \Delta_p \to 0 \). Given that, in our units, \( F(p) \) varies on the scale of unity, it is sufficient to take \( \Delta_p \ll 1 \). In practice, we take \( \Delta_p = 0.1 \). We also monitor the position, adding \( p \) to the position every time the particle takes a step, starting at momentum \( p \).

To simulate paths of the Langevin dynamics, we used both a straightforward Euler-Maruyama integration of the Langevin equation and one based on the random-walk picture above. Both methods give equivalent results.

2. The Bessel excursion

An efficient way to generate the set of Bessel excursions is to generate all discrete Brownian excursions with equal probability and then to weight each excursion by its appropriate weight. The task of generating the ensemble of Brownian excursions is, at first sight, a nontrivial problem since the constraint that the random walk not cross the origin is nonlocal in the individual left/right steps constituting the walk. However, it is easily accomplished using Callan’s proof [81] of the famed Chung-Feller theorem [82]. Callan gives an explicit mapping of any \( N \) left-step, \( N \) right-step random walks to a unique walk that does not cross the origin. This mapping maps exactly \( N + 1 \) of the \( \binom{2N}{N} \) \( N \) left, \( N \) right walks to each of the nonzero crossing walks [of which there are \( \binom{2N}{N} / (N + 1) \)]. Since every Brownian excursion of \( N + 2 \) steps consists of an initial right step, a non-zero-crossing \( N \) left-step, \( N \) right-step walk and a final left step back to the origin, Callan’s mapping allows us to generate an \( N + 2 \)-step Brownian excursion by generating a random \( N \) left-step, \( N \) right-step walk, applying the map, and prepending and postpending the appropriate steps. In this way, every excursion is generated with exactly the same probability.

Callan’s mapping is as follows. For a given \( N \) left-step, \( N \) right-step walk, one accepts it as is if it does not cross the origin. Otherwise, one finds the leftmost point reached by the walk (i.e., the minimal \( x \) of the path), which we denote \( x_n \), where \( n \) is the number of steps to reach this point. If this point is visited more than once, then one takes the first visit. One then constructs a new walk based on the original walk. We divide the original walk into two segments, the first \( n - 1 \) steps and the remaining part. The first part of the new walk consists of the second segment, shifted by \( -x_n \) such that the new walk starts on the origin. In the original walk, the step before reaching \( x_n \) is obviously a left step. One appends this left step at the end of the walk under construction and then attaches the first segment (of the original walk). This new walk clearly does not cross the origin, while preserving the number of left and right steps; hence, it starts and ends on the origin.

To account for the bias, i.e., to switch from Brownian to Bessel excursions, it is sufficient to weight each excursion by the weight factor \( F \):

\[
F = \prod_{i=1}^{N} f_i;
\]

\[
f_i = \begin{cases} 1 + \frac{F(p_i) \Delta_p}{2D} & \text{Step } i \text{ is a right step}. \\ 1 - \frac{F(p_i) \Delta_p}{2D} & \text{Step } i \text{ is a left step}. \end{cases}
\]

This works well as long as \( D \) is not too small. For very small \( D \), the convergence is quite slow since small areas have anomalously large weight unless \( N \) is very large. As long as one is interested in the large-\( N \) behavior, better convergence is achieved by taking \( F(p) = \theta(p - 1) / p \).

APPENDIX E: AREAL DISTRIBUTION OF THE BESSEL MEANDER

Our goal is to find the conditional PDF \( p_M(\chi | r^\tau) \), where \( \chi = \int_0^\tau p(t) \, dt \) is the area under the Bessel meander. The starting point for the calculation of \( p_M(\chi | r^\tau) \) is a modification of Eq. (20). As for the area under the Bessel excursion, we consider only the positive meander where \( 0 < \chi < \infty \), and later we use the symmetry of the process to find the areas under both positive and negative meanders. Since, contrary to the condition on the excursion, the meander is not bound to return to the origin, we now keep the end point free and integrate the propagator \( \hat{G}_C(s, p|p) \) over all possible values of \( p \). Hence,

\[
\hat{p}_M(s| r^\tau) = \lim_{\epsilon \to 0} \int_0^\infty \hat{G}_C(s, p|s) \, dp,
\]

with \( \hat{p}_M(s| r^\tau) = \int_0^\infty p_M(\chi | r^\tau) \exp(-sD\chi^\tau) \, d\chi^\tau \).

We expand the propagator \( \hat{G}_C(s, p|p) \) in a complete orthonormal basis, using the same approach as in the main text, while accounting for the boundary conditions of the meander. We thus rewrite Eq. (20) for the case of the meander and integrate over \( p \),

\[
\int_0^\infty \hat{G}_C(s, p|p) \, dp = e^{1/2D}s^{1/2}\int_0^\infty dp \, p^{-1/2D} \sum_k f_k(s^{1/3}e) \times f_k(s^{1/3}p) e^{-DE_i r^\tau},
\]

where the functions \( f_k(\cdot) \), as before, are the solutions of the time-independent equation (27) and \( E_k = \lambda_k^{1/2} \). Changing variables to \( \tilde{p} = s^{1/3}p \), and using the small \( p \) behavior, Eq. (28), which gives \( f_k(s^{1/3}e) \sim \tilde{d}_k s^{1/3}/16D e^{1/2D} \), we find

\[
\int_0^\infty \hat{G}_C(s, p|p) \, dp = e^{2s} \sum_k \tilde{d}_k e^{-D_0^{1/3} \lambda_k r^\tau} \int_0^\infty f_k(\tilde{p}) / \tilde{p}^{1/2D} \, d\tilde{p},
\]

for the numerator in Eq. (E1). We now define a new \( k \)-dependent coefficient.
\[ a_k = \int_0^\infty \frac{f_k(p)}{p^{1/2}} \, dp. \quad (E4) \]

The \( a_k \)'s, similar to the \( \tilde{d}_k \)'s, are evaluated from a numerically exact solution of Eq. (27). Rewriting in terms of the new coefficients,

\[ \int_0^\infty \tilde{G}_c(s, p|\epsilon) \, dp = e^{2\alpha \epsilon} \sum_k \tilde{d}_k a_k e^{-D_k \lambda_k \tau}. \quad (E5) \]

For the denominator in Eq. (E1), we expand the propagator in the orthonormal basis for the \( s = 0 \) case,

\[ \tilde{G}_c(s = 0, p|\epsilon) = \sum_k \left( \frac{e}{p} \right)^{1/2D} \phi_k^0(e) \phi_k^0(p) e^{-DE_k \tau}. \quad (E6) \]

Integrating over \( p \), as required for the meander, we get

\[ \int_0^\infty dp \tilde{G}_c(s = 0, p|\epsilon) = \int_0^\infty dp \sum_k \left( \frac{e}{p} \right)^{1/2D} \phi_k^0(e) \phi_k^0(p) e^{-DE_k \tau}. \quad (E7) \]

Plugging in \( \phi_k^0(e) \) and \( \phi_k^0(p) \) using Eq. (32), with \( B_k = \sqrt{\pi k/L} \), we find

\[ \int_0^\infty dp \tilde{G}_c(s = 0, p|\epsilon) = e^{1/2+1/2D} \frac{\pi}{L} \int_0^\infty dp \sum_k p^{-1/2D} J_{\alpha}(ke) \sqrt{\alpha} J_{\alpha}(kp) e^{-DE_k \tau}. \quad (E8) \]

Transforming the sum into an integral over \( k \) : \( \sum_k \to \frac{1}{\pi} \int_0^{\infty} dk \), and using the small \( p \) behavior of the Bessel function for \( J_{\alpha}(ke) \sim \frac{(ke)^\alpha}{\pi^{1/2} \Gamma(1+\alpha)} \), we get

\[ \int_0^\infty dp \tilde{G}_c(s = 0, p|\epsilon) = \frac{e^{2\alpha}}{2^{1/2} \Gamma(1+\alpha)} \int_0^\infty dp \, p^{1/2} \int_0^\infty dk k^{1+\alpha} J_{\alpha}(kp) e^{-DE_k \tau}. \quad (E9) \]

with \( E_k = k^2 \) in this continuum limit. To solve Eq. (E9), define \( q = kp \), and so

\[ \int_0^\infty dp \tilde{G}_c(s = 0, p|\epsilon) = \frac{e^{2\alpha}}{2^{1/2} \Gamma(1+\alpha)} \int_0^\infty dq \int_0^\infty dk \, q^{1/2-1/2D} J_{\alpha}(q) e^{-DE_k \tau}. \quad (E10) \]

Then, inserting \( x = k^2 \), we get

\[ \int_0^\infty \tilde{G}_c(s = 0, p|\epsilon) \, dp = \frac{e^{2\alpha}}{2^{1/2} \Gamma(1+\alpha)} \int_0^\infty dx x^{1/2-1} e^{-D \tau} \int_0^\infty dq q^{1/2-1/2D} J_{\alpha}(q). \quad (E11) \]

The two integrals on the right-hand side of Eq. (E11) can be evaluated analytically, and the final result for the denominator of Eq. (E1) is

\[ \int_0^\infty \tilde{G}_c(s = 0, p|\epsilon) \, dp = \frac{e^{2\alpha}}{2^{1/2} \Gamma(1+\alpha)} (D \tau)^{-\alpha}. \quad (E12) \]

Plugging Eqs. (E3), (E4), and (E12) into Eq. (E1), we reach our first main result for the area under the Bessel meander,

\[ \hat{p}_M(s|\tau^*) = 2^{2\alpha} \Gamma(1+\alpha) |s(D \tau^*)^{3/2}| \sum_k \tilde{d}_k a_k e^{-D_k \lambda_k \tau^*}. \quad (E13) \]

The inverse of \( \hat{p}_M(s|\tau^*) \), i.e., \( sD \to \chi^* \), is

\[ p_M(\chi^*|\tau^*) = 2^{2\alpha} \Gamma(1+\alpha) D^{3/2} |(\tau^*)^{3/2}| \sum_n \tilde{c}_n \tilde{d}_n \frac{(-c_k\chi^*)^{n-(2n/3)+1}}{\Gamma(-n+2n/3) n!}. \quad (E14) \]

with \( c_k = D^{1/3} \lambda_k \tau^* \). And finally, summing over \( n \) using Maple,
This is the main result of this appendix. In Fig. 14, we plot this areal distribution, comparing it with a histogram obtained from finite time simulations.

**APPENDIX F: THE TIME INTERVAL STRADDLING**

The velocity process \( v(t') \) restricted to the zero-free interval containing a fixed observation time \( t \) is called the excursion process straddling \( t \), and the portion of it up to \( t \) is the meandering process ending at \( t \). For Brownian and Bessel processes, these random paths were the subject of intense mathematical investigation; see, e.g., Refs. [17,31,83] and references therein. The investigation of the duration of the excursion straddling \( t \), the time interval for the meander ending at \( t \), and statistical properties of the path, e.g., its maximal height, have attracted mathematical attention since they reveal deep and beautiful properties of Brownian motion in dimension \( d \) (as is well known, the Bessel process is constructed from the radius of a \( d \)-dimensional Brownian motion). One aspect of the problem is the quantification of the properties of the set of points on the time axis on which the zero-crossing take place. Clearly, the number of zero crossings for a Brownian path starting at the origin within the time interval \((0, t)\) is infinite because of the continuous nature of the path (see more details below). One might naively expect that the time between points of zero crossing approaches zero since a finite measurement time divided by infinity is zero. That, in the context of our paper, might be true when \( \tau \) is finite. Then, we can expect to find a zero crossing in the close vicinity of the measurement time \( t \) (i.e., at a distance of the order of \( \tau \) from \( t \)). However, when \( 0 < \alpha < 1 \), the situation is more subtle. In this case, the points on the time axis will be clustered, and while their number is infinite, visualizing these dots with a simulation, we will observe a fractal dust.

In this appendix, we investigate the statistics of duration of the excursion straddling time \( t \), deriving some known results along the way. Our approach is based on renewal theory, using methods given by Godreche and Luck [47]. The previous mathematical approaches are based on direct analysis of Brownian and Bessel motion, while we use the \( \epsilon \) trick introduced in the main text. As mentioned, we replace the continuous regularized Bessel process or Brownian path with a noncontinuous one, with jumps of size \( \epsilon \) after zero crossing, and at the end, \( \epsilon \) is taken to zero.

Consider the PDF of waiting times \( g(\tau) \) with \( 0 < \tau < \infty \), with the long-time limiting behavior

\[
g(\tau) \sim \tau^{1/\alpha} \epsilon^{-(1+\alpha)}. \tag{F1}\]

While \( \alpha = (1 + D)/2D \) for the problem in the main manuscript (hence, \( \alpha \geq 1/2 \)), here we will assume that \( 0 < \alpha < 1 \). The case \( \alpha = 1/2 \) is the Brownian case since a particle starting on \( \epsilon \) will return for the first time to the origin according to the well-known law \( g(\tau) \propto \tau^{-3/2} \). In our physical problem, the \( \epsilon \) trick means that \( g(\tau) \) is \( \epsilon \) dependent, Eq. (59). As in the main text, we denote \( \{t_1, t_2, \ldots, t_k, \ldots\} \) as times of the renewal events (i.e., in our problem, zero crossing events). This is a finite set of times if \( \epsilon \) is finite. The waiting times \( \{\tau_1, \tau_2, \ldots\} \) are independent, identically distributed random variables because of the Markovian property of the underlying paths, and \( t_1 = \tau_1, t_k = \sum_{j=1}^{k-1} \tau_j \). We call \( t \) the observation time, and the number of renewals in \((0, t)\) is denoted by \( n \); thus, by definition, \( t_n < t < t_{n+1} \) (this is obviously true as long as \( \epsilon \) is finite). The time \( B = t - t_n \) is called the backward recurrence time [47], the time \( F = t_{n+1} - t \) is the forward recurrence time, and \( \Delta = t_{n+1} - t_n \) is the time interval straddling \( t \). In our problem, the backward recurrence time is the duration of the meandering process, which starts at \( t_n \) and ends at \( t \). For Brownian motion and Bessel processes, where \( t_n \) and \( t_{n+1} \) are the zero hitting times straddling \( t \), the statistics of \( F, B \), and \( \Delta \) are nontrivial.

We now obtain the PDF of \( \Delta \), denoted \( d_\delta(\Delta) \), using the methods in Ref. [47]. From the constraint \( t_n < t < t_{n+1} \), we have

\[
d_\delta(\Delta) = \sum_{n=0}^{\infty} [\delta(\Delta - (t_{n+1} - t_n))]I(t_n < t < t_{n+1})). \tag{F2}\]

where \( \delta[\cdot] \) is the Dirac \( \delta \) function; \( I(t_k < t < t_{n+1}) = 1 \) if the condition in the parentheses is true, and otherwise, it is zero. We use the double Laplace transform

\[
d_\epsilon(\alpha) = \int_{0}^{\infty} \int_{0}^{\infty} e^{-st} e^{-\alpha \Delta} d_\delta(\Delta) d\alpha d\Delta. \tag{F3}\]
and compute
\[
\left\langle \int_0^\infty \int_0^\infty e^{-st}e^{-u\Delta}\delta(\Delta-(t_{n+1}-t_n))f(t_n < t < t_{n+1})dtd\Delta \rightangle
=\frac{1}{s}\left((e^{-t\tau_0}-e^{-t\tau_{n+1}})e^{-u(t(t_{n+1}-t_n))}\right).
\]  
\hfill (F4)

Since the waiting times \(\{\tau_i\}\) are mutually independent, identically distributed random variables, we find, using \(t_n = \sum_{i=1}^n \tau_i\),
\[
(e^{-t\tau}) = \tilde{g}(s),
\]  
\hfill (F5)

where \(\tilde{g}(s) = \int_0^\infty e^{-st}g(t)dt\) is the Laplace transform of \(g(t)\). A similar result holds for the other expressions in Eq. (F4). It is then easy to find
\[
d_s(u) = \sum_{n=0}^\infty \frac{\tilde{g}(u) - \tilde{g}(u + s)}{s},
\]  
\hfill (F6)

and summing the geometric series, we get
\[
d_s(u) = \frac{\tilde{g}(u) - \tilde{g}(u + s)}{s[1 - \tilde{g}(s)]},
\]  
\hfill (F7)

When \(g(t) = \exp(-t)\), we find, in the limit of large \(t\), \(d_s(\Delta) \sim \Delta\exp(-\Delta)\), so the minimum of this PDF is at \(\Delta = 0\) and \(\Delta = \infty\), which is expected. For our more interesting case, \(\tilde{g}(u) \sim 1 - g^*/\Gamma(-\alpha)|u^\alpha\) for small \(u\); hence, we get, in the small \(u\) and \(s\) limits, with an arbitrary ratio between them,
\[
d_s(u) \sim \frac{(u + s)^\alpha - u^\alpha}{s^{1 + \alpha}}.
\]  
\hfill (F8)

This is a satisfying result, implying that the solution does not depend on \(g^*\), and hence it does not depend on \(\epsilon\). Thus, the artificial cutoff \(\epsilon\), which was introduced merely as a mathematical tool, does not alter our final formulas. Let \(\Delta = \Delta/t\) and denote by \(\tilde{d}(\Delta)\) its PDF in the long-time limit. With a useful inversion formula, given in the Appendix of Ref. [47], we invert Eq. (F8) to the time domain and find
\[
\tilde{d}(\Delta) = \frac{\sin \pi \alpha}{\pi (\Delta)^{1+\alpha}}[1 - (1 - \Delta)^\alpha I(0 < \Delta < 1)],
\]  
\hfill (F9)

which exhibits a discontinuity of its derivative at \(\Delta = 1\). Thus, the PDF of \(0 < \Delta < \infty\) has a cusp at \(t\), which allows for the identification of the measurement time \(t\) from a histogram of \(\Delta\). For Brownian motion, \(\alpha = 1/2\),
\[
\tilde{d}(\Delta) = \begin{cases} 
\frac{1}{2}\frac{1}{(\Delta)^{1/2}} & \Delta < 1 \\
\frac{1}{2}\frac{1}{(\Delta)^{1/2}} & \Delta > 1.
\end{cases}
\]  
\hfill (F10)

Instead of fixing \(t\), we may draw \(t_c > 0\) from an exponential PDF, \(R\exp(-Rt_c)\), in such a way that the mean of \(t_c\), \(\langle t_c \rangle = 1/R\), is very large; thus, the number of renewals in the time interval \([0, 1/R]\) is large. Similar to the previous case, we define \(\Delta_c = t_{n+1} - t_n\), which is called the duration of the interval straddling an independent exponential time, and \(t_n < t_c < t_{n+1}\). This case was treated rigorously by Bertoin et al. [31] for the Bessel process. Using the renewal theory approach, we now obtain their main result on the statistics of \(\Delta_c\), with a few hand-waving arguments. From the definition of the Laplace transform, \(\tilde{f}(s) = \int_0^\infty f(t_c)e^{-st}dt_c\), we see that, for a function \(f(t_c)\) depending on a random variable \(t_c\), the latter being exponentially distributed with mean equal to \(1/R\), the averaged function is the Laplace transform of \(f(t_c)\) evaluated at \(s = 1/R\) followed by multiplication with \(R\), namely, \(\langle f(t_c) \rangle = R\tilde{f}(R)\). Inserting \(s = R\) in Eq. (F7) followed by multiplication with \(R\), we get the Laplace transform of the PDF of \(\Delta_c\),
\[
\langle \exp(-u\Delta_c) \rangle = \langle [u + R]^{\alpha} - u^\alpha \rangle/R^\alpha,
\]  
\hfill (F11)

where we used the small \(R\) and \(u\) limit and, as usual, \(\langle \exp(-u\Delta_c) \rangle = \int_0^\infty e^{-u\Delta}d_s(\Delta_c)d\Delta\), where \(d_s(\Delta_c)\) is the PDF of \(\Delta_c\). This is the known result for the Bessel process when scaled properly, i.e., \(R = 1\) in Ref. [31], where the inverse Laplace transform of Eq. (F11) is also given, thus providing an explicit formula for the PDF \(d_s(\Delta_c)\).

To see the connection between renewal theory and statistics of the duration of Bessel excursion straddling an exponential time, notice that in Ref. [31] a Bessel process in dimension \(0 < d < 2\) is considered with the relation \(d = 2(1 - \alpha)\). In our work, we consider motion in a logarithmic potential in one dimension, which is easily mapped onto a Bessel process in dimension \(d = 1 - 1/D\). Thus, the exponent \(\alpha\) in Ref. [31] is the same as ours since, as we have shown, \(\alpha = (1 + D)/(2D)\), Eq. (30). Notice that for optical lattices \(D > 0\); hence, this system is a physical example for a regularized Bessel process in dimension \(-\infty < d < 1\). Finally, note that the original Bessel process considers the distance \(|r|\) from the origin of a Brownian motion in \(d\) dimensions, this being non-negative, \(|r| \geq 0\). Hence, the Bessel process does not exhibit zero crossings, so the points on the time axis are zero hitting points, not zero crossing points. This is a minor technical issue because of the symmetry of the binding effective potential discussed in this paper (i.e., negative and positive excursions in the logarithmic potential are statistically identical).
Similarly, we can obtain the limiting distributions of $B$ and $F$ using renewal theory [47]. The PDF of $\tilde{B} = B/t$ (here, $t$ is fixed) is

$$\text{PDF}(\tilde{B}) \sim \frac{\sin \pi \alpha}{\pi} (1 - \tilde{B})^{-\alpha} \quad 0 < \tilde{B} < 1, \quad (F12)$$

and the PDF of $\tilde{F} = F/t$, originally derived by Dynkin [84], is

$$\text{PDF}(\tilde{F}) \sim \frac{1}{\pi} \frac{1}{(\tilde{F})^{\alpha}(1 + \tilde{F})} \quad 0 < \tilde{F} < \infty. \quad (F13)$$

Inserting $\alpha = 1/2$, we get the result for $\tilde{B}$, for zero crossing of Brownian motion, obtained by Chung [see Eq. (2.22) there]. One can also quantify other aspects of zero crossings. For example, it is well known that the average number of renewals follows $\langle n \rangle \sim t^{1/3} [1 - (\alpha = 1/2)]$, and hence, using Eq. (59), $\langle n \rangle \propto t^{1/3}/e$, which diverges in the continuum limit $e \to 0$, as it should. The distribution of $n/\langle n \rangle$ is the well-known Mittag-Leffler distribution, with index $\alpha$ and unit mean.

To conclude, we see that the backward and forward recurrence times scale linearly with $t$, and they exhibit nontrivial behavior, which can be obtained either from analysis of Brownian or Bessel processes or by using renewal theory with the $e$ trick. The latter is a very simple approach, which requires some basic results in renewal theory. Importantly, given the tools of renewal theory, the approach, which requires some basic results in renewal theory with the analysis of Brownian or Bessel processes or by using recurrence times scale linearly with $t$.

$$F(\tau) \sim F(\tau) \sim \frac{1}{\pi} \frac{1}{(\tilde{F})^{\alpha}(1 + \tilde{F})} \quad 0 < \tilde{F} < \infty. \quad (F13)$$

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$\tilde{F} \sim F(\tau) \sim \frac{1}{\pi} \frac{1}{(\tilde{F})^{\alpha}(1 + \tilde{F})} \quad 0 < \tilde{F} < \infty. \quad (F13)$


[34] For atoms in molasses with a $J_y = 1/2 \rightarrow J_y = 3/2$ Zeeman substructure in a linear laser configuration, $c = 12.3$ [6]. References [5,19,35] give $c = 22$. As pointed out in Ref. [6], different notations are used in the literature.


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[70] L. F. Richardson, Atmospheric Diffusion Shown on a Distance Neighbour Graph, Proc. R. Soc. A 110, 709 (1926).