A Monte Carlo Method for 3D Radiative Transfer Equations with Multifractional Singular Kernels

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Abstract

We propose in this work a Monte Carlo method for three dimensional scalar radiative transfer equations with non-integrable, space-dependent scattering kernels. Such kernels typically account for long-range statistical features, and arise for instance in the context of wave propagation in turbulent atmosphere, geophysics, and medical imaging in the peaked-forward regime. In contrast to the classical case where the scattering cross section is integrable, which results in a non-zero mean free time, the latter here vanishes. This creates numerical difficulties as standard Monte Carlo methods based on a naive regularization exhibit large jump intensities and an increased computational cost. We propose a method inspired by the finance literature based on a small jumps - large jumps decomposition, allowing us to treat the small jumps efficiently and reduce the computational burden. We demonstrate the performance of the approach with numerical simulations and provide a complete error analysis. The multifractional terminology refers to the fact that the high frequency contribution of the scattering operator is a fractional Laplace-Beltrami operator on the unit sphere with space-dependent index.

Key words. radiative transfer, singular scattering kernels, Monte Carlo method, wave propagation, random media, long-range correlations.

1 Introduction

Radiative transfer models have been used for more than a century to describe wave energy propagation through complex/random media [32, 10], as well as neutron transport [40, 51], heat transfer [54], and are still an active area of research in astrophysics, geophysics, and optical tomography [39, 43, 44, 45] for instance. In this work, we propose a new Monte Carlo (MC) method to simulate the following radiative transfer equation (RTE)

\[
\begin{aligned}
\frac{\partial}{\partial t} u + \hat{k} \cdot \nabla_x u &= Q u, \\
u(t = 0, x, \hat{k}) &= u_0(x, \hat{k}),
\end{aligned}
\]

(1)

where \(S^2\) denotes the unit sphere in \(\mathbb{R}^3\), and \(u\) is the wave energy density in the context of wave propagation or a particle distribution function in the context of neutronics. The scattering operator \(Q\) has the standard form

\[
(Qu)(x, \hat{k}) = \lambda(x) \int_{S^2} \Phi(x, |\hat{\beta} - \hat{k}|)(u(x, \hat{\beta}) - u(x, \hat{k}))\sigma(d\hat{\beta}),
\]

(2)

for \(\sigma(d\hat{\beta})\) the surface measure on \(S^2\), \(\Phi\) the scattering kernel, and \(\lambda > 0\) a function modeling the support of the scattering process. Regions where \(\lambda(x) = 0\) are homogeneous and \(u\) undergoes free transport. MC methods have long be used for the resolution of (1), see e.g. [36, 51]. The originality and difficulty in our work lies in the fact that we consider situations where the mean free time \(t_0\) associated with \(Q\) vanishes in the scattering regions, that is

\[
\frac{1}{t_0(x)} = \lambda(x) \int_{S^2} \Phi(x, |\hat{\beta} - \hat{k}|)\sigma(d\hat{\beta}) = +\infty, \quad \text{where} \quad \lambda(x) > 0,
\]

(3)

and as a consequence the standard MC representations of \(u\) do not apply. Such a scenario arises for instance in the context of highly peaked-forward light scattering in biological tissues and in turbulent atmosphere, or more
generally in the context of wave propagation in random media with long-range correlations that we describe below. In this paper we write $\Phi$ as

$$
\Phi(x, |\hat{p} - \hat{k}|) := \frac{a(|\hat{p} - \hat{k}|)}{|\hat{p} - \hat{k}|^{2+\alpha(x)}} = \frac{1}{21+\alpha(x)/2} \rho(x, \hat{k} \cdot \hat{p}), \quad \text{with} \quad \rho(x, s) := \frac{a(\sqrt{2(1-s)}}{(1-s)^{1+\alpha(x)/2}} \quad s \in [-1, 1). (4)
$$

Above, $\alpha : \mathbb{R}^3 \rightarrow [0, 2)$ accounts for the slow variations of scattering across the ambient space, and $a$ is a smooth bounded function characterizing some statistical properties of the medium and such that $a(0) > 0$. Practical examples are given further. A direct calculation shows that (4) holds when $\alpha \in (1, 2)$. Also, the integral in (4) has to be understood in the principal value sense when $\alpha \in [1, 2)$, see [23]. The multifractional terminology that we use is motivated by the fact that the unbounded operator $Q$ can be expressed as a (multi)-fractional Laplace-Beltrami operator $(-\Delta_{\mathbb{R}^2})^{\alpha(x)/2}$ on the unit sphere up to a bounded operator w.r.t. the $\hat{k}$ variable [22, 23].

We would like to emphasize that we focus in this work on kernels of the form (4) for simplicity of the exposition, and that our method applies, after proper decomposition (see [23]), to more general kernels that behave like (4) at the singularity.

The RTE can be derived from high frequency wave propagation in random media, see e.g. [49]. In such a context, the velocity field $c(x)$ has the form

$$
\frac{1}{c^2(x)} = \frac{1}{c_0^2} \left(1 + \sqrt{\eta} V_0 \left(x, \frac{x}{\eta} \right)\right), \quad x \in \mathbb{R}^3, \quad \eta \ll 1,
$$

where $c_0$ is the background velocity (that we set to one in the sequel for simplicity), $V_0$ is a mean zero random field modeling fluctuations around the background, and $\eta$ is the correlation length of the random medium, assumed to be small after proper rescaling. The first variable in $V_0$ represents the slow variations of the random perturbations, while the second one corresponds to their high frequency oscillations. The latter are responsible for the strong interaction between the wave and the medium over sufficient distances. The scattering kernel $\Phi$ is related to the correlation function of $V_0$, and assuming $V_0$ is stationary (in the statistical sense) with respect to the fast variable, a kernel of the form (4) can be obtained from random fields such that

$$
\mathbb{E}[V_0(x, x') V_0(y, y')] = \sqrt{\lambda(x)\lambda(y)} \int_{\mathbb{R}^3} \frac{a(|p|)}{|p|^{1+2(\alpha(x)/2)}} e^{ip(x' - y')} dp, \quad (5)
$$

with $\alpha$ ranging from 0 to 2. Denoting by $R(x)$ the expectation in (5) with $y = x$, $y' = x' + x/\eta$, one can show that $R$ behaves like $|x|^{\alpha(x)-2}$ for $|x| \gg 1$, and is therefore not integrable. This is how random fields with long-range correlations are defined, as opposed to random fields with short-range correlations that exhibit an integrable correlation function. This approach is of practical interest in biomedical imaging as media with long-range correlations are able to reproduce experimentally observed power-law attenuations associated with effective fractional wave equations [20, 25]. The value of the exponents is related to the rate of decay of the correlation function $R$, and depends on the nature of the imaged tissues as reported in [14, 26, 27]. Variations of this exponent can then be used for diagnosis purposes [38, 47].

In Figure [1] we provide examples of such 2D random fields. The top-left picture represents a random medium with short-range correlations (with a standard Gaussian covariance kernel), while the top-right picture illustrates a random medium with long-range correlations with $\alpha \equiv 1$. Because of the singularity at $p = 0$, one can observe significantly larger statistical patterns than in the short-range case. In the bottom two pictures, we highlight the roles of $\lambda$ and $\alpha$: $\lambda$ characterizes scattering regions, and $\alpha$ defines the correlation structure. In the inner circle of the bottom-left picture we have $\alpha \equiv 0.1$, which tends to create shorter range fluctuations than in the outside where $\alpha \equiv 1$. In the bottom-right picture, we have a three-layer model for $\alpha$ in which the inner band exhibits smaller statistical patterns than the outer ones. This type of model is used for modeling non-Kolmogorov atmospheric turbulences, while standard atmospheric turbulence is modeled with the so-called Kolmogorov power spectrum

$$
\Phi(|k|) \propto \frac{a(|k|)}{|k|^{11/3}},
$$

for $|k|$ in the inertial range of turbulence. This corresponds to the case $\alpha = 5/3$. This case is not always valid in experiments as reported in [11, 52, 55], and the statistics of atmospheric turbulence have been shown to vary with altitude. Models have been derived for instance [see 35 for a review] by considering three ranges (0-2km, 2-8km, and above 8km) with distinct power laws (see Figure [10] for an illustration).
Figure 1: Realizations of Gaussian random fields. The upper-left picture represents a field with short-range correlations, while the upper-right picture depicts a field with long-range correlations with $\lambda = 1$ and $\alpha = 1$. In the lower-left picture, we have $\lambda = 1$ and $\alpha = 1$. In the lower-right, we have $\lambda = 1$ and $\alpha = 1$.

In the context of biological tissues, the following the Gegenbauer scattering kernel $\rho_G$ and Henyey-Greenstein (HG) kernel $\rho_{HG}$ are commonly used in the peaked-forward regime [29, 48]:

$$
\rho_G(x, s) := \frac{\alpha g (1 + g^2 - 2gs)^{-1/2}}{2\pi ((1 - g)^{-\alpha} - (1 + g)^{-\alpha})}, \quad \rho_{HG}(x, s) := \frac{1 - g^2}{4\pi (1 + g^2 - 2gs)^{3/2}}.
$$

The parameter $g \in (-1, 1)$ is called the anisotropy factor, and $\rho_{HG}$ is obtained by setting $\alpha = 1$ in $\rho_G$. The case $g = 0$ corresponds to isotropic energy transfer over the unit sphere, $g < 0$ to dominant transfer in the backward direction, and $g > 0$ to forward energy transfer. The peaked forward regime is obtained in the limit $g \to 1$, for which

$$
\frac{1}{(1 - g)^{\alpha}} \rho_G(x, \hat{k}, \hat{p}) \sim g^{\alpha} \frac{2\pi (2 - 2\hat{k} \cdot \hat{p})^{1+\alpha/2}}{2\pi (\hat{k} - \hat{p})^{2+\alpha}}.
$$

The case $\alpha = 1$ for the HG kernel is widely used in photon scattering in biological tissues [13, 21, 31]. A typical realization of the corresponding random field in 2D as $g \to 1$ is depicted in the top-right panel of Figure 1.

There exist a variety of methods for the resolution of (1) that handle the singular nature of the HG kernel, see e.g. [18, 19, 33, 34, 37]. They are based on finite differences type discretizations, projections over appropriate bases w.r.t. the $\hat{k}$ variable, and approximations of the kernel. Here we propose an alternative approach to handle singular scattering kernel (4) that is based on a MC method. The latter are popular choices for the simulation of the RTE when the kernel is smooth, see e.g. [36, 41, 42, 46, 51], essentially for their adaptability to a wide range of configurations and their simplicity of implementation. A downside is their slow convergence rate, and there is a vast literature on variance reduction techniques for acceleration. In this work, we focus on the design of an efficient MC method and postpone any variance reduction considerations to future works.

Our approach is based on an adaptation of a method proposed by Asmussen-Cohen-Rosiński [3, 11] (ACR) for the simulation of Lévy processes with infinite jump intensity. It relies on a small jumps/large jumps decomposition of the corresponding infinitesimal generator. The main idea is to approximate the generator of the small-jump part, which possesses the infinite intensity due to the singularity of kernel, by a Laplace-Beltrami operator (with respect to the angular variables) on the unit sphere $S^2$. This requires us to simulate paths of a jump-diffusion process over the unit sphere. For this purpose, we use the characterization of Brownian motion on the unit sphere given in [5] based on a standard stochastic differential equation (SDE) in $\mathbb{R}^3$ that is suitable for space-dependent kernels. This situation is hence more involved than the 2D case we investigated in [24]
where the small jumps part can be approximated by Brownian motion on the unit circle for which analytical expressions are available. Note that, as shown in [24], neglecting small jumps altogether in order to use standard MC methods leads to large errors, and reducing those comes at significantly increased computational cost.

Denoting by $\hat{\mu}(u)$ the estimator produced by our MC method for some observable $\mu(u)$ built on the solution $u$ to (1), we provide an error estimate of the form

$$P\left(|\hat{\mu}(u) - \mu(u)| > E_1 + E_2 + E_3\right) \ll 1$$

as a theoretical support of our method. Above, $E_1$, $E_2$, and $E_3$ are small terms characterizing the various approximation errors from the original model: the Laplace-Beltrami (i.e. small jumps) approximation, the discretization error of the diffusion process over the unit sphere, and the MC error. Note that the method we propose here applies directly to the stationary version of (1)

$$\dot{k} \cdot \nabla_x u - \mathcal{Q}u = u_0, \quad (x, \hat{k}) \in \mathbb{R}^3 \times \mathbb{S}^2,$$

with source term $u_0$, through the relation

$$u(x, \hat{k}) := \int_0^\infty u(t, x, \hat{k})dt.$$

The paper is organized as follows. In Section 2, we introduce probabilistic representations for (1) and its approximation based on the ACR method. In Section 3, we describe our MC method, state the main theoretical result regarding the overall approximation error, and detail the simulation algorithms. Section 4 is dedicated to the validation of the method using semi-analytical solutions. Numerical illustrations are given in Section 5, where we investigate the role of the strength $\alpha$ of the singularity, both when constant or space-dependent in the case of non-Kolmogorov turbulence, and compare with solutions for the HG kernel. Section 6 is devoted to the proofs of our main results and we recall in an Appendix the stochastic collocation method.

The numerical simulations are performed using the Julia programming language (v1.6.5) on a NVIDIA Quadro RTX 6000 GPU driven by a 24 Intel Xeon Silver 2.20GHz CPUs station. The codes have been implemented using the CUDA.jl library [8-9].

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## 2 Probabilistic representations and approximation

### 2.1 Representation for (1)

The starting point is the following standard probabilistic interpretation to (1):

$$u(t, x, k) = E_{x, k}[u_0(D(t))] := E[u_0(D(t)) | D(0) = (x, \hat{k})],$$

where $D = (X, K)$ is a Markov process on $\mathbb{R}^3 \times \mathbb{S}^2$ with infinitesimal generator

$$\mathcal{L}f(x, \hat{k}) := -\hat{k} \cdot \nabla_x f(x, \hat{k}) + \frac{\lambda(x)}{2^{1+\alpha(x)/2}} \int_{\mathbb{S}^2} \rho(x, \hat{p} \cdot \hat{k})(f(x, \hat{p}) - f(x, \hat{k}))\sigma(d\hat{p}).$$

A path, or a realization, of the Markov process $D$ is often referred to as a particle trajectory. The $X$ component of $D$ represents the position of a particle, and the component $K$ its direction. The generator $\mathcal{L}$ comprises two terms, the transport part describing free propagation of the particle, and the scattering operator (often referred to as the jump part in the probabilistic literature) describing the evolution of its direction. The jump component exhibits a non-integrable singularity leading to a infinite jump intensity and a vanishing mean free time as expressed in (3).

Note that when $\lambda$ and $\alpha$ are constant, it is shown in [23] that the solution $u$ is unique and infinitely differentiable in all variables for $t > 0$ for any square integrable initial condition. When $\lambda$ and $\alpha$ are infinitely differentiable with bounded derivatives at all orders, this result remains valid and we will assume throughout this work that $u$ is smooth. The same applies to the function $u_0$ defined further in Proposition 2.1.

In order to adapt the ACR method, we introduce the following small region over which the singularity of the kernel $\rho$ in (4) is not integrable, resulting in an unbounded infinitesimal generator $\mathcal{L}$:

$$S^\varepsilon_\hat{k} = S^\varepsilon(x|\hat{k}) := \{\hat{p} \in \mathbb{S}^2 : 1 - \hat{p} \cdot \hat{k} < \varepsilon\} \quad \varepsilon \in (0, 1).$$  

4
We can now decompose the jump part of the generator $\mathcal{L}$ into two components

$$\mathcal{L}f(x, \dot{k}) = -\dot{k} \cdot \nabla_x f(x, \dot{k}) + \mathcal{L}_{\infty}^e f(x, \dot{k}) + \mathcal{L}_d^f f(x, \dot{k})$$

$$:= -\dot{k} \cdot \nabla_x f(x, \dot{k}) + \frac{\lambda(x)}{2^{1+\alpha(x)/2}} \left( \int_{S^2_\infty} + \int_{S^2_\infty} \right) \rho(x, \dot{p} \cdot \dot{k})(f(x, \dot{p}) - f(x, \dot{k})) \sigma(d\dot{p}),$$

where $S^2_\infty = (S^2)^c$ is the complementary set of region $\mathcal{S}$ over the unit sphere. The part of the operator involving $S^2_\infty$ (with no singularity) is the infinitesimal generator of a standard jump Markov process. Regarding $S^2_\infty$ (with the singularity), the resulting approximation of this singular part by a Laplace-Beltrami operator $\Delta_{S^2}$ over the unit sphere $S^2$. We will use the notation $r_\varepsilon = \sqrt{1 - (1 - \varepsilon)^2/(2 - \varepsilon)}$ in what follows, and set in the rest of the paper $0 < \varepsilon \leq \varepsilon_0 < 1$ and $0 \leq \alpha_m \leq \alpha(x) \leq \alpha_M < 2$.

**Proposition 2.1** Let $u$ be the solution to (1) and $u_\varepsilon$ be the solution to

$$\begin{align*}
\partial_t u_\varepsilon + \dot{k} \cdot \nabla_x u_\varepsilon &= \sigma_\varepsilon^2(x) \Delta_{S^2} u_\varepsilon + \frac{\lambda(x)}{2^{1+\alpha(x)/2}} \int_{S^2_\infty} \rho(x, \dot{p} \cdot \dot{k})(u_\varepsilon(\dot{p}) - u_\varepsilon(\dot{k})) \sigma(d\dot{p}), \\
u_\varepsilon(0, x, \dot{k}) &= u_0(x, \dot{k}),
\end{align*}$$

for $(t, x, \dot{k}) \in (0, \infty) \times \mathbb{R}^3 \times S^2$, where

$$\sigma_\varepsilon^2(x) := \frac{2^{1-\alpha(x)} \alpha(0) \pi \lambda(x)}{2 - \alpha(x)} r_\varepsilon^{2-\alpha(x)}.$$

Assuming $a'(0) = 0$, for any $T > 0$, we have

$$\sup_{t \in [0, T]} \|u(t, \cdot, \cdot) - u_\varepsilon(t, \cdot, \cdot)\|_{L^2(\mathbb{R}^3 \times S^2)} \leq \varepsilon^{2-(\alpha_M/2)} \sqrt{2T} E(u)$$

where $E(u)$ is defined in [34].

The proof of Proposition 2.1 is postponed to Section 6.1. The term $E(u)$ is independent of $\varepsilon$ and depends on derivatives of $u$ w.r.t. $\dot{k}$ up to order 4. Note that the error is of order $\varepsilon^{3-(\alpha_M/2)}/(2 - \alpha_M)$ when $a'(0) \neq 0$ yielding a less accurate approximation than for $a'(0) = 0$. The difference comes from a truncated expansion along the sphere covering providing an extra order in $\varepsilon$ assuming $a'(0) = 0$. This later assumption holds throughout the remaining of the paper. Based on (11), we then devise a MC method for (9) instead of (1). The advantage in using (9) is the fact that the angular diffusion term $\sigma_\varepsilon^2(x) \Delta_{S^2}$ is the generator of a Markov process that can be easily simulated. Indeed, for $W$ a standard 3D Brownian motion on $\mathbb{R}^3$ and $\times$ the cross product in $\mathbb{R}^3$, it is shown in [32] that the process $B$ solving the SDE

$$dB = B \times dW - B dt, \quad B(0) \in S^2,$$

has generator $\frac{1}{2} \Delta_{S^2}$. A simple adaptation then gives the desired diffusion coefficient. Since the error is of order $\varepsilon^{2-(\alpha_M/2)}$, it is always smaller than $\varepsilon$, and can be adjusted to obtain a desired accuracy. Note also that $\sigma_\varepsilon(x)$ increases as $\alpha(x)$ gets to 2, and diffusion on the sphere eventually becomes the dominant dynamics.

### 2.2 Representation for [4]

We interpret [4] as the forward Kolmogorov equation of an appropriate Markov process, and as a consequence focus on forward MC methods, see e.g. [36] for terminology. Backward equations are simulated in a similar manner, and can be combined with forward methods for variance reduction techniques [71 40 51].

The Markov process we consider for this approach is defined by

$$D_\varepsilon(t) := \sum_{n \geq 0} 1_{[T_n, T_{n+1})}(t) \psi_n^{Z_n}(t - T_n) \quad t \geq 0,$$

where (in the remaining of the paper we extensively make use of the notation $z = (x, \dot{k})$):

1. The flow $\psi_n^{Z_n}$ is the unique strong solution to the SDE

$$\begin{align*}
dX_n^Z(t) &= K_n^Z(t) dt, \\
dK_n^Z(t) &= \sqrt{2} \sigma_\varepsilon(X_n^Z(t)) K_n^Z(t) \times dW_n(t) - 2 \sigma_\varepsilon^2(X_n^Z(t)) K_n^Z(t) dt,
\end{align*}$$

where $\psi_n^Z(0) = z$, $\times$ is the cross product in $\mathbb{R}^3$, $(W_n)_n$ is a sequence of independent standard Brownian motions on $\mathbb{R}^3$, and $\sigma_\varepsilon$ is defined by (10).
2. The jump times \((T_n)_n\) are distributed according to
\[
\mathbb{P}(T_{n+1} - T_n > t \mid D_z(T_n) = z, (\psi_n^\varepsilon(s))_{s \in [0,t]}) = \exp \left( -\int_0^t \Lambda_z(\psi_n^\varepsilon(s)) \, ds \right), \quad \forall n \geq 0,
\]
with \(T_0 = 0\), and for \(\rho\) given by (4),
\[
\Lambda_z(z) := \frac{\lambda(x)}{21 + \alpha(z)^2} \int_{S^2} \rho(x, \hat{\rho} \cdot \hat{k}) \sigma(\hat{\rho}).
\]  \hfill (14)

3. The jumps \((Z_n)_n\) describe a Markov chain with transition probability
\[
\mathbb{P}(Z_{n+1} \in dy \otimes \sigma(\hat{d}\rho) \mid Z_n, T_{n+1} - T_n) = \Pi_z(z_{n+1}, dz),
\]  \hfill (15)
where \(z_{n+1} := \psi_n^\varepsilon(T_{n+1} - T_n)\), and
\[
\Pi_z(z, dz) := \pi_z(z, \hat{\rho}) \sigma(\hat{\rho}) \delta_z(dy),
\]  \hfill (16)
with density
\[
\pi_z(z, \hat{\rho}) := \frac{\rho(x, \hat{\rho} \cdot \hat{k})}{\int_{S^2} \rho(x, \hat{\rho} \cdot \hat{k}) \sigma(\hat{\rho})} \mathbf{1}_{S^2}(\hat{\rho}), \quad z = (x, \hat{k}),
\]  \hfill (17)
which is supported over \(S^2\). The above Dirac mass \(\delta_z(dy) := \delta(x-y)dy\) translates the fact that the jumps only hold w.r.t. the \(\hat{k}\) variable.

Let us note that the above family of standard Brownian motions \((W_n(t))_{t \in [0, T_n + 1 - T_n]}\) can be defined as
\[
W_n(t) = W(t + T_n) - W(T_n), \quad t \in [0, T_n + 1 - T_n],
\]
for any \(n\), where \(W\) is a single standard Brownian motion on \(\mathbb{R}^3\). We have then the following probabilistic representation for the solution to (6).

**Proposition 2.2** The Markov process \(D_z\) defined in (11) has for infinitesimal generator
\[
A_z g := \hat{k} \cdot \nabla_x g(z) + \sigma^2(z) \Delta g(z) + \Lambda_z(z) \int_{S^2} \pi_z(z, \hat{\rho})(g(x, \hat{\rho}) - g(x, \hat{k})) \sigma(\hat{\rho}),
\]
and we have
\[
\mathbb{P}_{\mu_0}(D_z(t) \in dx \otimes \sigma(d\hat{k})) = \frac{1}{u_0} u_z(t, x, \hat{k}) \, dx \, \sigma(d\hat{k}),
\]  \hfill (18)
where
\[
\mu_0(dx, d\hat{k}) := \mathbb{P}(D_z(0) \in dx \otimes \sigma(d\hat{k})) = \frac{u_0(x, \hat{k})}{u_0} \, dx \, \sigma(d\hat{k}) \quad \text{with} \quad \hat{u}_0 := \int_{\mathbb{R}^3 \times S^2} u_0(x, \hat{k}) \, dx \, \sigma(d\hat{k}).
\]  \hfill (19)

The terminology *forward* comes from the fact that the particles are emitted at random points at time \(t = 0\) (through \(\mu_0\)) and propagate towards the observation position \(z = (x, \hat{k})\). The proof of Proposition 2.2 is provided in Section 6.2. Let us illustrate two aspects of the representation (18). In order to obtain an estimation of \(u_z(t, x, \hat{k})\) at the point \(z = (x, \hat{k})\), we calculate the probability
\[
\frac{u_0}{|B(z, r)|} \mathbb{P}_{\mu_0}(D_z(t) \in B(z, r)) \simeq u_z(t, x, \hat{k}),
\]  \hfill (20)
where \(B(z, r) \subset \mathbb{R}^3 \times S^2\) stands for the open ball centered at \(z = (x, \hat{k})\) with radius \(r \ll 1\). If we are only interested in e.g. the energy density at point \(x\), we estimate
\[
\frac{\tilde{u}_0}{|B(x, r)|} \mathbb{P}_{\mu_0}(D_z(t) \in B(x, r) \times S^2) = \frac{\tilde{u}_0}{|B(x, r)|} \mathbb{P}_{\mu_0}(D_{1,x}(t) \in B(x, r)) \simeq \int_{S^2} u_z(t, x, \hat{k}) \sigma(d\hat{k}),
\]
where \(B(x, r) \subset \mathbb{R}^3\) stands for the open ball centered at \(x\) with radius \(r \ll 1\), and \(D_{1,x}\) is the \(x\) component of \(D_z\).
3 Monte Carlo Method

Based on the previous probabilistic representation of (9), solving (9) requires the generation of random paths of the stochastic process \(D_z\). For any measurable bounded functions \(f\), the convergence of the estimator
\[
\mu_N(t, f) := \frac{1}{N} \sum_{j=1}^N f(D^j_z(t)) \xrightarrow{N \to \infty} \int f(x, \dot{k}) u_z(t, x, \dot{k}) \, dx \sigma(d\dot{k}) \quad \mathbb{P}_{\mu_0} \text{-- almost surely},
\]
is guaranteed by the strong law of large numbers. Above \((D^j_z)_j\) is a sample of \(D_z\). We detail next how to treat efficiently the diffusion and jump components of the process \(D_z\).

3.1 The jump part

Since the process \(D_z\) is inhomogeneous, i.e. \(\Lambda_z\) and \(\Pi_z\) both depend on \(z = (x, \dot{k})\), we use the so-called thinning method, also referred to as the fictitious shocks method [36]. It is based on a acceptance/rejection step and consists in simulating at first more jumps (or shocks) than necessary. In a second step, some of the jumps are rejected according to an appropriate probability distribution in order to recover the original dynamics. Assume \(0 < \alpha_m \leq \alpha(x) \leq \alpha_M < 2\). A direct calculation shows that
\[
\Lambda_z(z) \leq \frac{2\pi \sup \lambda \sup a}{21+\alpha(z)/2} \int_{-1}^{1-\varepsilon} \frac{dt}{(1-t)^{1+\alpha(z)/2}} = \frac{4\pi \sup \lambda \sup a}{\alpha(x)^{21+\alpha(z)/2} \varepsilon^{\alpha(z)/2}} \leq \frac{2\pi \sup \lambda \sup a}{\alpha_m \varepsilon^{\alpha_M/2}} := \breve{\Lambda}_z. \tag{21}
\]
The fictitious jump times are then drawn as
\[
\bar{T}_n := \sum_{j=1}^n \xi_j, \quad \text{and} \quad \bar{T}_0 = 0,
\]
where the \((\xi_j)_j\) are i.i.d. exponentially distributed random variables with parameter \(\breve{\Lambda}_z\).

The thinning method consists in the following acceptance/rejection step. At a jump time \(\bar{T}_n\) and current position \(z_n\), we draw a jump \(z\) according to the probability distribution \(\Pi_z\) given by \([13]\). This jump is accepted with probability \(p(z_n) = \Lambda_z(z_n)/\Lambda_z\). Otherwise, the process \(D_z\) continues to diffuse starting from \(z_n\), and \(\bar{T}_n\) is not considered as a true jump time for \(D_z\). Practically, we can define the state as
\[
\hat{Z}_n := z \mathbb{1}_{(U_n \leq p(z_n))} + z_n \mathbb{1}_{(U_n > p(z_n))}
\]
at each fictitious jump times \(\bar{T}_n\). Above, \(U_n\) is a random variable uniformly distributed over \([0, 1]\) and all the \(U_n\)’s are independent.

3.2 The diffusion part

The diffusion part between two jumps satisfies the linear SDE \([13]\), and is simulated using the following Euler-Maruyama type scheme
\[
(S_{n,m}) : \begin{cases}
X_{n,m+1} = X_{n,m} + h_{n,m} \dot{K}_{n,m} \\
K_{n,m+1} = \dot{K}_{n,m} - 2h_{n,m} \sigma^2(X_{n,m}) \dot{K}_{n,m} + \sqrt{2h_{n,m} \sigma(X_{n,m})} \dot{K}_{n,m} \times W_{n,m} \\
\end{cases} \tag{22}
\]
where the \((W_{n,m})_{m,n}\) are i.i.d. mean-zero Gaussian random vectors with identity covariance matrix. Note that the above scheme does not conserve the Euclidean norm with respect to the angular variable, and as a consequence the evolution of \((K_{n,m})_{n,m}\) does not remain on the unit sphere over the iterations. This motivates the definition of \(K_{n,m}\). We have nevertheless \(\mathbb{E}[|K_{n,m}|^2] = 1\), for all \(n\) and \(m\), and Theorem 3.1 below guarantees that the distribution of \(K_{n,m}\) provides a converging approximation of the true statistics. The stepsizes \(h_{n,m}\) are determined from a fixed stepsize \(h\) as follows. Since our convergence theorem further is stated at a fixed time \(T\), we include \(t = T\) in the discretization grid for simplicity. Let then \(N_T\) such that \(\bar{T}_{N_T} \leq T < \bar{T}_{N_T+1}\). When \(n \neq N_T\), let \(d_t = \bar{T}_{n+1} - \bar{T}_n\). With \(m_n = [d_t/h]\) (\([\cdot]\) the integer part), we set, for \(n \neq N_T\),
\[
h_{n,m} = \begin{cases}
h & \text{for } m = 0, \ldots, m_n - 1 \\
d_t - m_n h & \text{for } m = m_n,
\end{cases}
\]
and if \(m_n = 0\), we set \(h_{n,0} = d_t\). In the rest of the paper, the grid is denoted by \((T_{n,m})_m\), where \(T_{n,m+1} = T_{n,m} + h_{n,m}\) for \(T_{n,0} = \bar{T}_n, n \geq 0\) and \(m = 0, \ldots, m_n - 1\). When \(n = N_T\), we divide the interval \([\bar{T}_{N_T}, \bar{T}_{N_T+1}]\) similarly into subintervals of length \(h_{N_T,m}\) at most \(h\) (we suppose there are \(m_{N_T}\) of those) and such that \(\bar{T}_{N_T,m} = T\) for one \(m\) in \(0, \ldots, m_{N_T}\).
3.3 The overall discretized process and convergence

For any $t \geq 0$, the approximate version of the process $D_z$, denoted $D_{h,\varepsilon}$, is defined by:

$$D_{h,\varepsilon}(t) = (X_{h,\varepsilon}(t), K_{h,\varepsilon}(t)) := \sum_{n=0}^{\infty} \sum_{m=0}^{m_{n}} 1_{[T_{n,m}, T_{n,m}+\varepsilon_{n,m}]}(t)Z_{n,m},$$

where

1. For any $m \in \{0, \ldots, m_{n}\}$,

$$Z_{n,m+1} = S_{n,m}(Z_{n,m}),$$

where $Z_{n,0} = \tilde{Z}_{n}$ for the $(\tilde{Z}_{n})_{n}$ defined below, and where $S_{n,m}(Z_{n,m}) = (X_{n,m+1}, \dot{K}_{n,m+1})$ is given by the scheme (22) with initial condition $Z_{n,m} = (X_{n,m}, \dot{K}_{n,m}).$

2. The sequence $(\tilde{Z}_{n})_{n}$, is defined by

$$\tilde{Z}_{n+1} = \tilde{z} 1_{\{U_{n} \leq p(Z_{n,m_n+1})\}} + Z_{n,m_n+1} 1_{\{U_{n} > p(Z_{n,m_n+1})\}} \quad n \geq 0,$$

where $\tilde{z}$ is drawn according to the probability measure $\Pi_{\varepsilon}(Z_{n,m_n+1}, dz)$ defined by (16).

Below, $b$ is the backward solution to (1) with terminal condition $b(T, x, \dot{k}) = f(x, \dot{k})$, see (41). Our convergence result is then the following (we set $h$ such that $4h \sup_x \sigma^2(x) \leq 1$ to simplify some expressions):

**Theorem 3.1** Consider

$$\mu_{N,h,\varepsilon}(t, f) = \frac{1}{N} \sum_{j=1}^{N} f(D_{h,\varepsilon}^{j}), \quad \mu(t, f) = \int f(x, \dot{k}) u(t, x, \dot{k}) dx \sigma(d\dot{k}),$$

where $(D_{h,\varepsilon}^{j})$ is a sample of $D_{h,\varepsilon}$. For any $T > 0$, $\eta > 0$ and any smooth bounded function $f$ on $\mathbb{R}^3 \times \mathbb{S}^2$, we have

$$\limsup_{N \to \infty} \mathbb{E} \left( |\mu_{N,h,\varepsilon}(T, f) - \mu(T, f)| \right) \leq \frac{\eta \Sigma_{h,\varepsilon}}{\sqrt{N}} + \varepsilon^{2-(\alpha_M/2)} \sqrt{\text{Var}(f)} + h F_0(u, b, f) + h F_1(b) \leq \text{erfc}(\eta/\sqrt{2}),$$

(23)

where

$$\Sigma_{h,\varepsilon} = \sqrt{\text{Var}(f(D_{h,\varepsilon}(T)))} \leq \sup |f|.$$

The functions $F_0$ and $F_1$ are explicit and independent of $\varepsilon$ and $h$, and are defined in the proof of the theorem in Section 6.3.

Theorem 3.1 is proved in Section 6.3. In (23), there are three terms that quantify the approximation error of our estimator $\mu_{N,h,\varepsilon}(t, f)$: one of order $\varepsilon^{2-(\alpha_M/2)}$ due to the approximation of $u$ by $u_{\varepsilon}$ (the smaller the $\alpha_M$, i.e. the less singular the kernel is, the smaller the error), one of order $h$ due to the numerical approximation of the diffusion over the unit sphere, and one due to the MC approximation with the standard $1/\sqrt{N}$ convergence rate. Note that the discretization error of the diffusion process is only of order $h$ and not of order the standard $\sqrt{h}$. The reason is that we are only interested in the convergence of Monte Carlo estimators, allowing us to consider this discretization error in the weak sense [3]. However, a weak second-order Runge-Kutta method can be considered to provide an error in $h^2$ instead of $h$ for the Euler scheme [12]. Modifications of the SDE [13] can also be considered to provide weak higher-order scheme [1]. The main goal of this paper being to present a methodology to capture efficiently the behavior induced by the singularity, we focus our attention on the error in $\varepsilon$, and do not present weak-higher order discretization schemes for the SDE. In this way, the Euler scheme is considered for simplicity in the proof of Theorem 3.1. For the numerical simulations of Sections 4 and 5 that illustrate the roles of $\varepsilon$ and $\alpha$ in the approximation, the parameter $h$ will be chosen proportionally to the shortest mean free time $\Lambda_{\varepsilon}^{-1}$, and small enough so that the approximation error w.r.t. $\varepsilon$ is dominant. $N$ will also be chosen large enough so that the error of approximation in $\varepsilon$ is dominant. The MC error is controlled by the standard deviation $\Sigma_{h,\varepsilon}$, and variance reduction techniques can be designed to reduce this term. When estimating the energy density over a given region $B$, as in (20), the number of particles $N$ needed to reach a given error threshold can be estimated as follows: the root mean square error of the MC estimator for $f = 1_B/|B|$ reads

$$\text{RMSE}_{h,\varepsilon} := \frac{\mathbb{E} \left[ (\mu_{N,h,\varepsilon}(T, f) - \mathbb{E}[f(D_{h,\varepsilon}(T))])^2 \right]^{1/2}}{\mathbb{E}[f(D_{h,\varepsilon}(T))]^{1/2}} = \sqrt{\frac{P_{h,\varepsilon}(1 - P_{h,\varepsilon})}{N |B|^2}} \leq \frac{1}{2\sqrt{N |B|}},$$

(24)
and the relative MC error is

$$\frac{\text{RMSE}_{h,\varepsilon}}{\mathbb{E}[f(D_{h,\varepsilon}(T))]} = \frac{1}{\sqrt{N}} \sqrt{\frac{1 - P_{h,\varepsilon}}{P_{h,\varepsilon}}} \leq \frac{1}{\sqrt{N}P_{h,\varepsilon}},$$

(25)

with

$$P_{h,\varepsilon} := \mathbb{P}_{\mu_0}(D_{h,\varepsilon}(T) \in B) \approx \frac{1}{\bar{u}_0} \int_B u_\varepsilon(t, x, k) dx \sigma(d\hat{k}).$$

Above, $\mu_0$ and $\bar{u}_0$ are given by (19). A RMSE lower than a threshold $c$ would then require

$$N \geq \frac{1}{4c^2|B|^2},$$

(26)

while a relative error would require

$$N \geq \frac{1}{c^2P_{h,\varepsilon}}.$$ 

(27)

If $B$ is a region centered around a point $(x_0, k_0)$, with a small volume (that is $P_{h,\varepsilon} \ll 1$ as for (20)), we would have

$$N \geq \frac{\bar{u}_0}{c^2|B|u_\varepsilon(t, x_0, k_0)}.$$ 

### 3.4 Algorithms

We discuss in this section practical aspects of the method. Before stating the algorithm itself, let us emphasize that a key point is to sample efficiently the jumps from $\Pi_e$ given by (16).

Let us fix the current state of the process $D_{h,\varepsilon}$ at a point $z = (x, k)$. In spherical coordinates, $\pi_\varepsilon$ defined in (17) is equivalent to a probability density function drawing a polar angle $\theta$ and an azimuthal angle $\varphi$. Here, the north pole of the spherical system is the current direction $\hat{k}$, and it is direct to see that the azimuthal angle $\varphi$ is uniformly distributed over $(0, 2\pi)$. We denote this by $\varphi \sim \mathcal{U}(0, 2\pi)$. For the polar angle, a change of variables leads to $\theta = \arccos(1 - \chi)$, where $\chi$ has probability density function

$$f_\chi(\chi|x) := \frac{a(\sqrt{2\chi})}{C_\chi \chi^{1+\alpha(x)/2}} 1_{(\varepsilon, 2)}(\chi),$$

and $C_\chi$ is a normalizing constant. Therefore, to draw a jump according to (17) starting from $\hat{k}$, we compute

$$\hat{\rho} = R(\theta, \varphi, \hat{k}) := \cos(\theta)\hat{k} + \sin(\theta)\left(I_3 + \sin(\varphi)Q(\hat{k}) + (1 - \cos^2(\varphi))Q^2(\hat{k})\right)\hat{k}^\perp,$$

(28)

where $\hat{k}^\perp$ is an orthonormal vector to $\hat{k}$, $I_3$ is the $3 \times 3$ identity matrix, and

$$Q(k) = \begin{pmatrix} 0 & -k_3 & k_2 \\ k_3 & 0 & -k_1 \\ -k_2 & k_1 & 0 \end{pmatrix}, \quad \text{where} \quad k = (k_1, k_2, k_3).$$

(29)

The transformation $R$ corresponds to a rotation from $\hat{k}$ to $\hat{\rho}$ with polar angle $\theta$ with respect to $\hat{k}$ and azimuthal angle $\varphi$ with respect to $\hat{k}^\perp$. Note that the choice of $\hat{k}^\perp$ is not important since $\varphi$ is uniformly distributed over $(0, 2\pi)$.

We notice that in the case of a constant function $a \equiv a_0$, one obtains a truncated Pareto distribution for $\chi$. The corresponding cumulative distribution function can be exactly inverted giving then a direct simulation method. In this case, the cumulative distribution function is given by, for $\chi \in (\varepsilon, 2)$,

$$F_\chi(\chi|x) = \frac{a_0}{C_\chi} \int_{\varepsilon}^{\chi} \frac{dv}{v^{1+\alpha(x)/2}} = \frac{1 - (\varepsilon/\chi)^{\alpha(x)/2}}{1 - (\varepsilon/2)^{\alpha(x)/2}},$$

The random variable $\chi$ can then be generated by

$$\chi = F_\chi^{-1}(U|x) = \varepsilon(1 - (1 - (\varepsilon/2)^{\alpha(x)/2}U)^{-2/\alpha(x)},$$

where $U$ is a random variable uniformly distributed over $(0, 1)$ ($U \sim \mathcal{U}(0, 1)$). In the case of a non constant function $a$, the main features of the density $f_\chi(\cdot|x)$ are similar to those of the truncated Pareto distribution, and
a stochastic collocation method can be considered to simulate $f_X(|x|)$. This method is described in Appendix A in our context. It is based on the simulation of the above truncated Pareto distribution and proves to be very effective.

The algorithm used to simulate a trajectory of $D_{h,\varepsilon}$ can be summarized in the following two procedures. The first one corresponds to the simulation of the diffusion process between two (fictitious) jumps, and we use the notation

$$S(Z, W, h) = \begin{cases} X + h K \\ K - 2h \sigma^2(X) K + \sqrt{2h} \sigma(X) W \end{cases}$$

with $Z = (X, K)$. Below, $N(0, I_3)$ stands for the three dimensional multivariate Normal distribution with identity covariance matrix.

### Algorithm 1: Diffusion

Input: current state of the particle $z = (x, \hat{k})$, duration of the diffusion $\delta t$

Output: state of the particle after the diffusion process

**Initialization:**

$\delta t \sim \mathcal{E}(\Lambda_{\varepsilon})$ // number of iterations

$Z \leftarrow z$ // initialization of the diffusion state

// Main loop of the diffusion

for $j \leftarrow 1$ to $n$

$W \sim N(0, I_3)$

$Z \leftarrow S(Z, W, h)$

$k \leftarrow K/|K|$ // Add a diffusion step with stepsize $h' \leq h$ to match the duration $\delta t$

$h' \leftarrow \delta t - nh$

$W \sim N(0, I_3)$

$Z \leftarrow S(Z, W, h')$

$k \leftarrow K/|K|$ // remaining diffusion step of duration $T - t$

return $Z$

### Algorithm 2: TrajectorySimulation

Input: Duration $T$ of the particle evolution

Output: state of the particle at time $T$

**Initialization:**

$Z \leftarrow (x, \hat{k}) \sim \mu_0$ // initialization of the particle state at random

$t \leftarrow 0$ // temporary time variable

$\delta t \sim \mathcal{E}(\Lambda_{\varepsilon})$ // first jump time

// main loop for the path evolution

while $t + \delta t < T$

$Z \leftarrow \text{Diffusion}(Z, \delta t)$

$U \sim \mathcal{U}(0, 1)$

if $U \leq p(Z)$ then

// the jump is accepted, $Z$ is transformed

$\chi \sim f_X(|x|)$

$\theta \leftarrow \arccos(1 - \chi)$

$\varphi \sim \mathcal{U}(0, 2\pi)$

$\hat{p} \leftarrow R(\theta, \varphi, \hat{k})$

$Z \leftarrow (x, \hat{p})$

$t \leftarrow t + \delta t$

$\delta t \sim \mathcal{E}(\Lambda_{\varepsilon})$

// remaining diffusion step of duration $T - t$

return $Z$

The rest of the paper is dedicated to numerical simulations and the proofs of our main results.
4 Validation

In this section, we first derive a semi-analytical solution to validate our method in the simplest situation where \( \alpha, a, \lambda \) are constant functions. We then highlight the crucial role of the small jumps correction for computational efficiency.

4.1 Semi-analytical solution

We set \( \lambda \equiv 1 \) and the RTE (1) reads

\[
\partial_t u + \hat{k} \cdot \nabla_x u = Qu
\]  

with scattering kernel

\[
Qf(\hat{k}) = a \int_{S^2} \frac{\sigma(d\hat{p})}{|\hat{p} - \hat{k}|^{2+\alpha}}(f(\hat{p}) - f(\hat{k})), \quad \hat{k} \in S^2.
\]

Using the Funk-Hekke formula \([50]\), this operator can be diagonalized in \( L^2(S^2) \) equipped with the inner product

\[
\langle f, g \rangle_{L^2(S^2)} = \int_{S^2} f(\hat{p})g(\hat{p})\sigma(d\hat{p}) = \int_0^\pi \int_0^{2\pi} f(\theta, \varphi)g(\theta, \varphi)\sin(\theta)d\theta d\varphi.
\]

The eigenvalues are given by

\[
\lambda_l = \frac{a \pi \Gamma(-\alpha/2)}{2^{\alpha} \Gamma(1+\alpha/2)} \left( \frac{\Gamma(l+1+\alpha/2)}{\Gamma(l+1-\alpha/2)} - \frac{\Gamma(1+\alpha/2)}{\Gamma(1-\alpha/2)} \right) \quad l \in \mathbb{N},
\]

and the eigenvectors are the spherical harmonics

\[
Y_{l,m}(\hat{k}) = Y_{l,m}(\theta, \varphi) := \sqrt{\frac{(2l+1)(l-m)!}{4\pi(l+m)!}} P_l^m(\cos(\theta))e^{im\varphi}, \quad (l, m) \in \mathbb{N} \times \{-l, \ldots, l\},
\]

where the \( P_l^m \) are the associated Legendre polynomials. In order to derive a semi-analytical solution, we Fourier transform \([30]\) w.r.t. \( x \), and introduce

\[
\tilde{u}(t, q, k) = \int_{\mathbb{R}^3} u(t, x, k)e^{-iq \cdot x} dx.
\]

Above, \( q = \tilde{q} := (0, 0, \xi) \) so that \( \tilde{u}(t, \xi, \hat{k}) = \tilde{u}(t, \tilde{q}, \hat{k}) \) solves

\[
\partial_t \tilde{u} + i \hat{k} \cdot \tilde{q} \tilde{u} = Qu.
\]  

Writing \( \hat{k} \) in spherical coordinates with \((0, 0, 1)\) as north-pole, this latter equation reads,

\[
\partial_t \tilde{u}(t, \xi, \theta, \varphi) = (Q - i\xi \cos(\theta))\tilde{u}(t, \xi, \theta, \varphi), \quad (t, \xi, \theta, \varphi) \in (0, \infty) \times \mathbb{R} \times (0, \pi) \times (0, 2\pi).
\]

We now decompose \( \tilde{u} \) on the basis of spherical harmonics

\[
\tilde{u}(t, \xi, \theta, \varphi) = \sum_{l=0}^{\infty} \sum_{m=-l}^{l} \tilde{u}_{l,m}(t, \xi)Y_{l,m}(\theta, \varphi),
\]

resulting in

\[
\frac{d}{dt} \tilde{u}_{l,m} = \lambda_l \tilde{u}_{l,m} - i\xi(d_{l,m}^+ \tilde{u}_{l+1,m} + d_{l,m}^- \tilde{u}_{l-1,m}) \quad \text{for} \quad l \geq 1, \quad \frac{d}{dt} \tilde{u}_{0,0} = -i\xi d_{0,0}^+ \tilde{u}_{1,0} \quad \text{for} \quad l = 0.
\]  

Above, we have used the fact that

\[
\langle Y_{m', \varphi, \cos(\theta)}Y_{l, \theta}, \rangle_{L^2(S^2)} = \begin{cases}
    d_{l,m}^+ := \sqrt{\frac{(l+m+1)(l-m+1)}{(2l+1)(2l+3)}} & \text{if} \quad m = m' \text{ and } l' - l = 1, \\
    d_{l,m}^- := \sqrt{\frac{(l+m)(l-m)}{(2l-1)(2l+1)}} & \text{if} \quad m = m' \text{ and } l' - l = -1, \\
    0 & \text{otherwise}
\end{cases}
\]
For computational purposes, we introduce a cutoff in the variable \( l \ (l \in \{0, \ldots, L\}) \), and consider a truncated version of (32) as the vector differential equation

\[
\frac{d}{dt} \hat{u}^L(t, \xi) = (D^L - i\xi A^L) \hat{u}^L(t, \xi), \quad \hat{u}^L(t, \xi) = (\hat{u}_{l_j+1}^L(t, \xi))_{l \in \{0, \ldots, L\}, j \in \{0, \ldots, 2l\}} \in \mathbb{C}^{(L+1)^2},
\]

where \( D^L \) and \( A^L \) are two \((L + 1)^2 \times (L + 1)^2\) matrices defined by

\[
\begin{align*}
D_{l_j+1, p_j+1}^L & := \lambda_l \quad \text{for} \quad l \in \{0, \ldots, L\}, \quad j \in \{0, \ldots, 2l\}, \\
A_{l_j+1, (l+1)(j+2)+1}^L & := d_{l_j-l-1}^L \quad \text{for} \quad l \in \{0, \ldots, L-1\}, \quad j \in \{0, \ldots, 2l\} \\
A_{l_j+1, (l-1)(j+2)+1}^L & := d_{l_j-l+1}^L \quad \text{for} \quad l \in \{1, \ldots, L\}, \quad j \in \{0, \ldots, 2(l-1)\}.
\end{align*}
\]

All other coefficients in both \( D^L \) and \( A^L \) are set to 0. Note that the indexing of the matrices starts at 0 for simplicity. The solution to (33) reads \( \hat{u}^L(t, \xi) = e^{(t-\xi A^L)}\hat{u}^L(0, \xi) \), where the matrix exponential is computed numerically. For our test case, we consider the following initial condition

\[
u(t = 0, x, k) = \frac{1}{\sqrt{2\pi}} e^{-|x|^2/2}, \quad 2 \cos^2(\theta/2) = \frac{1}{\sqrt{2\pi}} e^{-|x|^2/2} (2\sqrt{\pi}Y_{0,0}(\theta, \varphi) + 2\sqrt{\pi/3}Y_{0,1}(\theta, \varphi)),
\]

so that

\[
\hat{u}_{l_j+1}^L(t = 0, \xi) = \begin{cases} 
2\sqrt{\pi e^{-\xi^2/2}} & \text{for} \quad l = j = 0, \\
2\sqrt{\pi/3 e^{-\xi^2/2}} & \text{for} \quad l = j = 1, \\
0 & \text{otherwise}.
\end{cases}
\]

Finally, an approximation of \( \hat{u} \), solution to (31), is given by

\[
\hat{u}^L(t = 0, \xi, \theta, \varphi) = \sum_{l=0}^{L} \sum_{j=0}^{2l} e^{(D^L - i\xi A^L)t} \hat{u}^L(0, \xi)_{l_j+1} \ Y_{l, j-l}(\theta, \varphi).
\]

For numerical comparisons with our MC method, we introduce a discretization of the unit sphere \( S^2 \) via the polar and azimuthal angles \((\theta_m)_m\) and \((\varphi_m)_m\), with respective stepsize \( \Delta \theta \) and \( \Delta \varphi \). We then compare

\[
\hat{u}^L(t, \xi, \theta_m, \varphi_m') \simeq \frac{1}{\Delta \theta \Delta \varphi} \int_{\theta_m}^{\theta_{m+1}} \int_{\varphi_{m'}}^{\varphi_{m'+1}} \hat{u}^L(t, \xi, \theta, \varphi) \sin(\theta) d\theta d\varphi
\]

with its MC approximation

\[
\hat{u}^L_{N}(t, \xi, m, m') = \frac{1}{\Delta \theta \Delta \varphi N} \sum_{n=1}^{N} e^{-i\xi X_{\theta, \varphi}^n(t)} \ 1_{(\theta, \varphi) \in (\theta_m, \theta_{m+1}) \times (\varphi_{m'}, \varphi_{m'+1})}
\]

where \( \theta_{\theta, \varphi}^n \) and \( \varphi_{\theta, \varphi}^n \) are respectively the polar and azimuthal angles for \( K_{\theta, \varphi}^n \), and where \((D^L_{h, \varepsilon})_n = (X_{h, \varepsilon}^n, K_{h, \varepsilon}^n)_n\) is a sample of \( D_{h, \varepsilon} \) introduced in Section 3.

In the following numerical illustrations we consider \( \alpha = 0.002 \) in the RTE, and set \( \Delta \theta = \Delta \varphi = 0.05, \varepsilon = 0.1 \) and \( h = 0.5/\Lambda \simeq 12.6 \) for the approximation parameters. Note that these choices for \( \varepsilon \) and \( h \) are providing us with a good accuracy at a very low computational cost as we will see. Such values may have to be decreased in other setups and when considering different observables. For instance, in Section 5.2 where \( \alpha \) is varying, smaller values of \( \varepsilon \) and \( h \) are needed to capture correctly the solution.

Also, in the context of singular scattering kernels, the classical notion of scattering mean free time is not informative since it is equal to 0 (see (3)). Instead, we define a characteristic time using the inverse of the second eigenvalue of \( Q \), i.e. the first non zero eigenvalue, and set \( t_c = -\frac{1}{\lambda_1} \). We refer to Figure 2 for the evolution of \( t_c \) w.r.t. \( \alpha \).

In our setting, \( t_c \simeq 79.6 \) (for \( \alpha = 1 \)), which is about six times the stepsize \( h \) needed to capture the diffusive correction. Also, since \( \varepsilon \) is not too small, this correction plays a significant role in obtaining the correct dynamics.

In Figure 3, we compare, for \( \alpha = 1 \), the real and imaginary parts of the observable

\[
\hat{u}_1(t, \xi) := \int_0^\pi \int_0^{2\pi} \hat{u}^L(t, \xi, \theta, \varphi) \sin(\theta) d\theta d\varphi \quad \text{with} \quad 4\pi \Delta \theta \Delta \varphi \sum_{m,m'} \hat{u}^L_{N}(t, \xi, m, m'),
\]
Figure 2: Illustration of the evolution of $t_c$ w.r.t. $\alpha$.

Figure 3: Comparisons of the real and imaginary parts of $\hat{u}_1(t,\xi)$ for three observation times and for $\alpha = 1$. The grid in $\xi$ range from $-0.2$ to $0.2$ with 100 discretization points and we run $N = 2.4 \times 10^6$ particles.

Figure 4: Comparisons of the real and imaginary parts of $\hat{u}_2(t,\xi,\theta)$ for $\alpha = 1$, for three values of $\xi$ if $t = t_c$, and for $\xi = 0$ if $t = 10t_c$. We run $N = 2.4 \times 10^6$ particles for the top two pictures and $N = 24 \times 10^6$ for the third one.

for three values of $t$. In Figure 4 we compare the real and imaginary parts of

$$\hat{u}_2(t,\xi,\theta) := \int_0^{2\pi} \tilde{u}_L^L(t,\xi,\theta,\varphi)d\varphi \sin(\theta) \quad \text{with} \quad 4\pi \Delta \varphi \sum_{m'} \tilde{u}_N(t,\xi,m,m'),$$
for three values of $\xi$.

In Figure 5, we compare \( \hat{u}_3(t = 2t_c, \xi = 0.02, \theta) := \int_{0}^{2\pi} \tilde{u}_L^L(t = 2t_c, \xi = 0.02, \theta, \varphi) d\varphi \sin(\theta) \) with \( 4\pi \Delta \varphi \sum_{m'} \tilde{u}_N^L(t = 2t_c, \xi = 0.02, m') \), for three values of $\alpha$.

Figure 5: Comparisons of the real and imaginary parts of the observable $\hat{u}_3(t = 2t_c, \xi = 0.02, \theta)$ for three values of $\alpha$. We run $N = 24 \times 10^6$ particles.

In all these illustrations, and despite somewhat fairly large values for $\varepsilon$ and $h$, we observe a very good agreement between the Monte Carlo results and the semi-analytic calculations.

4.2 Role of the correction

In this section, we highlight the role of the correction provided by the diffusion over the unit sphere w.r.t. the $k$-variable. To this end, we compare the following observables obtained from the semi-analytic solution

\[
    u_4(x_3) = \int_{0}^{3t_c} \tilde{u}_L(t, x_\perp, x_3, \theta, \varphi) \sin(\theta) dt \, dx_\perp \, d\theta \, d\varphi, \quad x := (x_\perp, x_3) \in \mathbb{R}^2 \times \mathbb{R},
\]

with the ones obtained with our MC method, with and without this diffusive correction, and for various values of $\alpha$, $\varepsilon$ and $h$. The grid in $z$ range from $-300$ to $300$ with size $2^8$ and we run $N = 300 \times 10^6$ particles. According to (26) and (27), the number of samples $N$ is taken large enough so that the RMSE (24) of the MC estimation is of order $10^{-5}$ and the relative MC error (25) is of order $0.03\%$ where $u_4$ takes values of order as low as $10^{-3}$.

With this choice of $N$, we can focus our attention on the role played by $\varepsilon$ and $h$ in the approximation. In Figure 6, we represent the relative error $\text{Err}_{\varepsilon, \alpha}$ and running time of the MC method with and without a diffusive correction. The reference time in the right picture is the one of corrected method with $\varepsilon = 0.1$.

Figure 6: Illustration of the relative error $\text{Err}_{\varepsilon, \alpha}$ and running time of the MC method with and without a diffusive correction. The reference time in the right picture is the one of corrected method with $\varepsilon = 0.1$.

We represent the relative error

\[
    \text{Err}_{\varepsilon, \alpha} := \max_z \frac{|u_4(z) - u_{4,\text{MC}}(z)|}{\max_z u_4(z)},
\]

for various sizes of the cutoff $\varepsilon$, and where $u_{4,\text{MC}}(z)$ is the MC approximation to $u_4(z)$. The left picture illustrates the evolution of the relative error for various $\varepsilon$. The blue curve corresponds to the corrected MC
with $\varepsilon = 0.1$ (with still a fairly large stepsize $h = 0.5/\Lambda_\varepsilon$) providing at most a relative error slightly larger than $1\%$. The other curves correspond to the noncorrected MC method for several values of $\varepsilon$. The corrected MC consistently yields a better accuracy than the noncorrected version, and even in weakly singular cases where $\alpha$ is less than one, a very small value of $\varepsilon$ (red and green curves) is necessary to match the accuracy of the corrected method. The right picture illustrates the evolution of the relative running time of the noncorrected method w.r.t. the corrected one. For values of $\alpha$ less than 0.7 (weakly singular kernels), corrected and noncorrected methods have similar computational times for comparable accuracy, while in the case of singular kernels with $\alpha \geq 1$, the noncorrected methods yield a much larger cost and a much lower accuracy.

Figure 7: Illustration of the relative error $\text{Err}_{\varepsilon, \alpha}$ and running time of the (corrected) Monte Carlo method.

In Figure 8 we illustrate the precision and running time sensitivity of the (corrected) MC method w.r.t. the stepsize $h = h_0/\Lambda_\varepsilon$. As expected, we obtain a better precision for smaller stepsizes but at the price of a longer running time. These effects are amplified as $\alpha$ increases due to the increasing strength of the diffusion correction. In what follows, we select $h_0 = 0.3$ since this yields a relative error less than $1\%$ for a wide range of $\alpha$’s while not changing significantly the running time.

Figure 8: Illustration of the relative error $\text{Err}_{\varepsilon, \alpha}$ and running time of the (corrected) Monte Carlo method.

In Figure 8 we depict the precision and running time sensitivity w.r.t. the cutoff parameter $\varepsilon$, and observe the same phenomena as in the case of the stepsize $h$. The parameter $\varepsilon$ defines not only the accuracy of the diffusion correction, but also the average number of jumps, and as a consequence the running time increases as $\varepsilon$ decreases as in the case of the noncorrected Monte Carlo method.

5 Numerical illustrations

5.1 The role of $\alpha$

In this section, we highlight the effects of the kernel singularity on the energy density. We consider a constant $\alpha$, with $\alpha = 0.002$ in this section. Our setting is depicted in Figure 9. The spatial variable $x$ is decomposed into a main propagation axis $x_3$ and a transverse plane $x_\perp$, i.e. $x = (x_\perp, x_3) \in \mathbb{R}^2 \times \mathbb{R}$. The same notation holds for the direction variable $\hat{k} = (\hat{k}_\perp, \hat{k}_3) \in S^2$. We choose an initial condition for (1) of the form

$$u_0(x, \hat{k}) = \delta(x)\delta(\hat{k} - \hat{k}_0), \quad \hat{k}_0 = (0, 0, 1),$$
modeling a source located at $x = 0$ and embedded in the random medium, and emitting in the forward $x_3$-direction. We set a function $\lambda$ of the form $\lambda(x) = 1_{[-5,40]}(x_3)$, that defines a scattering layer between $x_3 = -5$ and $x_3 = 40$. In such a configuration, both transmitted and reflected quantities at $x_3 = 40$ and $x_3 = -5$ are of interest. With our particular choice for $k_0$, what is obtained at $x_3 = -5$ is purely due to backscattering.

In the following two subsections, the MC estimators for these quantities are given respectively by

$$
\hat{F}_{tr}(x_\perp) := \int_0^T dt \int_{S^2} \sigma(\hat{k}) \ u(t,x_\perp,x_3 = 40,\hat{k}) \quad \text{and} \quad \hat{F}_{ref}(x_\perp) := \int_0^T dt \int_{S^2} \sigma(\hat{k}) \ u(t,x_\perp,x_3 = -5,\hat{k}).
$$

The MC estimators for these quantities are given respectively by

$$
\hat{F}_{tr}(m,n) := \frac{1}{\Delta x_\perp N} \sum_{j=1}^N \mathbf{1}_{(X_{h,e}^j(\tau_i) \in \square_{mn}, X_{3,h,e}^j(\tau_i) > 40)} , \quad \hat{F}_{ref}(m,n) := \frac{1}{\Delta x_\perp N} \sum_{j=1}^N \mathbf{1}_{(X_{h,e}^j(\tau_i) \in \square_{mn}, X_{3,h,e}^j(\tau_i) < -5)}
$$

where

$$
\tau^j := \inf(t \in [0,T] : X_{3,h,e}^j(t) > 40 \quad \text{or} \quad X_{3,h,e}^j(t) < -5),
$$

is the first time the $j$-th particle escapes the slab. Note that once a particle escapes, it cannot reenter it since it propagates freely. Above, $(\square_{mn})_{m,n}$ is a uniform square grid of the traverse plane to the $x_3$-axis. All squares

\begin{figure}[h]
\centering
\includegraphics[width=0.5\textwidth]{figure9}
\caption{Illustration of the numerical setting.}
\end{figure}
in the grid have area $\Delta x_\perp$. Note that the grid can be different for the transmitted and reflected signals. We have considered for the transverse variable of the transmitted energy a uniform grid over a detector of size $[-10, 10] \times [-10, 10]$ centered around the $x_3$-axis, and over a detector of size $[-50, 50] \times [-50, 50]$ for the reflected energy. For both cases, we chose $128 \times 128$ grid points. The principle of these estimators is simply to count the number of particles that exit the slab before time $T$ and to record their position in the transverse plane.

In Figure 10 we illustrate the transmitted and reflected energy flux, for several values of $\alpha$. We represent the variations w.r.t. the first coordinate of $x_\perp = (x_1, x_2)$, and for two values of $x_2$.

![Figure 10: Illustration of the energy at the boundaries $F_{\text{tr}}^T(x_\perp)$ and $F_{\text{ref}}^T(x_\perp)$ w.r.t. $x_1$ with $x_2 = 0$ (solid line) and $x_2 = 1$ (dotted line) for the left picture and $x_2 = 10$ (dotted line) for the right one. The RMSEs (24) are less than $6.5 \times 10^{-4}$ (resp. $2.6 \times 10^{-5}$) on the left picture (resp. right picture), while the relative errors (25) are less than $0.6\%$ for the left picture (resp. $1.3\%$ for the right picture) for values of the observables as low as $10^{-3}$ (resp. $10^{-5}$).](image)

One can observe that at fixed times, the larger the $\alpha$, the more diffuse are the signals. Indeed, as $\alpha$ increases, the jump intensity $\Lambda_\varepsilon$ (in other words the number of scattering events) increases as well as the strength of the diffusive correction $\sigma_\varepsilon$ (see Figure 11).

![Figure 11: Illustrations of the evolution w.r.t. $\alpha$ of $\Lambda_\varepsilon$ defined by (14) and $\sigma_\varepsilon$ defined by (10). Here, $\alpha = 0.002$ and $\varepsilon = 0.01$.](image)

5.1.2 Time evolution of the exiting energy

Here, we are interested of the time evolution of the energy exiting the slab, and we define the (integrated) reflected and transmitted energy by

$$F_{\text{tr}}(t) := \int_{\mathbb{R}^2} dx_\perp \int_{S^2} \sigma(\hat{k}) u(t, x_\perp, x_3 = 40, \hat{k})$$

and

$$F_{\text{ref}}(t) := \int_{\mathbb{R}^2} dx_\perp \int_{S^2} \sigma(\hat{k}) u(t, x_\perp, x_3 = -5, \hat{k}).$$

The MC estimators for these two quantities are given by

$$\hat{F}_{\text{tr}}(n) := \frac{1}{dt N} \sum_{j=1}^{N} \mathbf{1}(\tau^j \in (t_n, t_{n+1}], \Gamma_{-h, \varepsilon}(\tau^j) > 40)$$

and

$$\hat{F}_{\text{ref}}(n) := \frac{1}{dt N} \sum_{j=1}^{N} \mathbf{1}(\tau^j \in (t_n, t_{n+1}], \Gamma_{-h, \varepsilon}(\tau^j) < -5).$$
Here, \((t_n)\) is a uniform grid of the time interval with stepsize \(dt\). For the transmitted signal, we have considered the time interval \([40, 45]\) with a stepsize \(dt = 0.02\), and have set \([0, 4t_c]\) with a stepsize \(dt = 0.4\) for the backscattered signal. Note that the time interval starts at 40 for the transmitted energy, which is the travel time of the wave (traveling at speed \(c_0 = 1\)) from the source to the plane \(x_3 = 40\). These estimators count the number of particles that exit the slab in the time interval \([t_n, t_{n+1}]\) at each side of the slab. In Figure 12 we illustrate the evolution of the transmitted and reflected energy, for several values of \(\alpha\).

![Figure 12: Time evolution of the energy at the boundaries \(F_{tr}(t)\) and \(F_{ref}(t)\). The RMSE (24) are less than \(8 \times 10^{-4}\) (resp. \(4 \times 10^{-5}\)) on the left picture (resp. right picture), while the relative errors (25) are less than 0.7\% for the left picture (resp. 0.5\% for the right picture) for values of the observables as low as \(10^{-3}\) (resp. \(10^{-4}\)).](image)

In the case of the transmitted signal (left), and for small values of \(\alpha\), we see the arrival of the coherent wave at the proper travel time followed by the coda. When \(\alpha\) increases, one notices the stronger impact of scattering and of the diffusive correction that smooths the signal out and damps its amplitude. For the largest \(\alpha\), we only observe a coda. Regarding the reflected signal (right), there is only a coda for all \(\alpha\) due our choice of \(\hat{k}_0\), and one can observe two stages in the dynamics: backscattering increases up to a time of order \(t_c\), about which exponentially decay due to the operator \(Q\) takes over.

### 5.1.3 Comparison with the Henyey-Greenstein scattering kernel

In this section, we compare the solutions to the RTE with Henyey-Greenstein scattering kernel \(6\) for an anisotropy factor \(g\) close to one with the solutions to \(9\) with singular kernel derived from \(7\), that is by setting \(a \equiv (1 - g)/(2\pi)\) and \(\alpha = 1\) in \(4\). Note that the value of the constant \(a\) changes with \(g\), and as a consequence \(\Lambda_\varepsilon, h,\) and \(\sigma_\varepsilon\) vary accordingly. To illustrate this approximation, we still consider the setting depicted in Figure 9 and the various observables introduced in the previous sections, but now at a time \(T = 300\).

We observe in Figure 13 the very good agreement between the two solutions. The reflected signal is well captured by our method despite fairly large values of \(\varepsilon\) and \(h\). Also, let us mention that the computational cost is decreasing as the anisotropic parameter \(g\) is getting close to 1, as the overall jump intensity decreases in this case in the highly peaked regime \(g \rightarrow 1\). Regarding the transmitted signal, \(\varepsilon\) (and then \(h\)) needs to be lowered for an accurate approximation, as explained at the beginning of Section 5.1.

The RTE with a Henyey-Greenstein scattering kernel is simulated with a standard MC method. Compared to our method, its computational costs to achieve RMSEs of order \(10^{-4}\) and \(10^{-5}\) for respectively the transmitted and reflected observables are the following:

| running time (s) | \(g = 0.97\) | \(g = 0.98\) | \(g = 0.99\) |
|------------------|--------------|--------------|--------------|
| HG kernel        | 8.3          | 7.6          | 6.0          |
| singular kernel, \(\varepsilon = 0.01\) | 13.9         | 6.7          | 2.0          |
| singular kernel, \(\varepsilon = 0.1\) | 2.5          | 1.5          | 0.7          |

Here, \(\varepsilon = 0.01\) is considered for the transmitted observables, while we set \(\varepsilon = 0.1\) for the reflected ones. According to this table, lower computational times are observed with our method for the three considered \(g\)’s compared to standard MC methods for the Henyey-Greenstein scattering kernel. Our MC method provides therefore an efficient tool to simulate an RTE with a Henyey-Greenstein kernel. For transmitted observables, \(g\) needs to be quite close to one to provide a significant advantage to our method.
Figure 13: Comparison of the observables obtained using the Henyey-Greenstein scattering kernel and our singular kernel with \( \alpha = 1, T = 300, \varepsilon = 0.01 \) for the transmitted observables (left panels), and \( \varepsilon = 0.1 \) for the reflected ones (right panels). The RMSEs (24) and relative errors (25) are similar to those of Figures 10 and 12.

5.2 Varying \( \alpha \) function

In this section, we investigate the influence of a varying \( \alpha \) function that characterizes the strength of the singularity. We consider two situations, one inspired from optical tomography, and the second one from wave propagation through atmospheric turbulence.

5.2.1 A two-stage model with a sphere

We keep the setting introduced in Section 5.1 and add a defect with a different value of \( \alpha \) to the setting. This defect is modeled by ball of radius 3 centered at the origin and where \( \alpha \) is equal to \( \alpha_1 \). We set \( \alpha \equiv 1 \) in the exterior of the ball, corresponding to the peak forward regime of the Henyey-Greenstein scattering kernel. See Figure 14. This situation models a biological tissue in which statistical properties are changing and define a region of interest for imaging.

![Figure 14](image)

Figure 14: Illustration of the setting with \( \lambda = 1_{\{x_3 \in (-5,40)\}} \) and \( \alpha(x) = \alpha_1 1_{x \in B} + 1 \cdot 1_{x \notin B} \) where \( B \) is a ball centered at 0 with radius 3.

We illustrate in Figure 15 the impact of the introduction of the defect on the observables introduced in Section 5.1. The impact is stronger on transmitted observables and quite significant, giving then the possibility to identify the defect with \( \alpha = \alpha_1 \) inside the scattering medium. Reflected quantities tend to be less sensitive to the presence of the defect since a fraction of the signal is backscattered before reaching it.
Figure 15: Illustrations the transmitted (left-hand-side) and reflected (right-hand-side) observables with \( T = 300 \), and \( \varepsilon = 0.01 \). For the top two pictures we set \( x_1 = 0 \) (solid lines) and \( x_1 = 1.5 \) for the top-left and \( x_1 = 0.5 \) for top-right picture (dotted lines).

5.2.2 Non-Kolmogorov turbulences

In this section, we keep once more the setting introduced in Section 5.1, with the difference that \( \alpha \) takes three different large values depending on the altitude parametrized by \( x_3 \), see Figure 16:

\[
\alpha(x_3) = \frac{5}{3} \cdot 1_{\{x_3 \leq 2\}} + \frac{4}{3} \cdot 1_{\{2 < x_3 \leq 8\}} + 1.9 \cdot 1_{\{8 < x_3\}}.
\]

The value \( 5/3 \) corresponds to standard Kolmogorov turbulences, while other values are associated with non-Kolmogorov turbulence models \([4, 52, 55]\). In these models, it is considered that for altitudes higher than 8km, the atmospheric turbulence yields larger statistical patterns (which tend to be created by singular kernels) than around the ground (0-2km). Hence, we set \( \alpha = 1.9 \) for altitudes greater than 8km. The function \( \alpha \) is no longer constant in these models, and for our illustrations we chose

\[
a(r) = 0.002 \cdot \exp\left(-r^2/(2 \times 0.8^2)\right).
\]

In Figure 17, one can notice that non-Kolmogorov turbulence yields quite different signals compared to Kolmogorov turbulence, in particular for reflected quantities. As we saw in Section 5.1, the higher the \( \alpha \), the more diffuse is the signal which then enhances reflected signals. This explains the increased reflections in the non-Kolmogorov case.

6 Proofs

This section is dedicated to the proof of Proposition 2.1, describing the approximation of the RTE (1) by (9) where the small jumps have been replaced by a diffusion term, Proposition 2.2, providing the probabilistic representation to (9), and Theorem 3.1, justifying the overall MC method involving a discretization scheme for the diffusion part.

6.1 Proof of Proposition 2.1

Let \( v_\varepsilon := u_\varepsilon - u \), so that \( v_\varepsilon(t = 0) = 0 \). We have

\[
\frac{d}{dt} \|v_\varepsilon(t)\|_{L^2(\mathbb{R}^3 \times S^2)}^2 = 2 \left< \partial_t v_\varepsilon(t), v_\varepsilon(t) \right>_{L^2(\mathbb{R}^3 \times S^2)}
\]

\[
= 2 < (\sigma_\varepsilon^2 \Delta_{\mathbb{S}^2} + L_\varepsilon^2) v_\varepsilon(t), v_\varepsilon(t) >_{L^2(\mathbb{R}^3 \times S^2)} + 2 < (\sigma_\varepsilon^2 \Delta_{\mathbb{S}^2} - L_\varepsilon^2) u, v_\varepsilon(t) >_{L^2(\mathbb{R}^3 \times S^2)}.
\]
Figure 16: Illustration of a three stages $\alpha$-profile for a non-Kolmogorov phase function.

Figure 17: Illustrations the transmitted (left-hand-side) and reflected (right-hand-side) observables with $T = 300$, and $\varepsilon = 0.01$. For the top two pictures we illustrate $x_2 = 0$ (solid lines) and $x_2 = 2$ for the top-left and $x_2 = 5$ for top-right picture (dot line).

Since $\Delta_{S^2}$ is a nonpositive operator, we have

$$< (\sigma_2^2 \Delta_{S^2} + \mathcal{L}_\varepsilon^c) v_c(t), v_c(t) >_{L^2(\mathbb{R}^3 \times S^2)} \leq -\frac{1}{2} \iiint_{\mathbb{R}^3 \times S^2} dx \sigma(d\tilde{p}) \sigma(d\tilde{k}) \frac{\lambda(x)}{2(1+\alpha(x)/2)} \int_{S^2} \rho(x, \tilde{p}, \tilde{k})(f(x, \tilde{p}) - f(x, \tilde{k}))^2 \leq 0.$$  

We then obtain

$$\sup_{t \in [0, T]} \|v_c(t)\|^2_{L^2(\mathbb{R}^3 \times S^2)} \leq 2 \int_0^T \|\sigma_2^2 \Delta_{S^2} - \mathcal{L}_\varepsilon^c\) u(t)\|^2_{L^2(\mathbb{R}^3 \times S^2)} dt,$$

which concludes the proof using the following lemma.

**Lemma 6.1** Let $0 < \varepsilon < \varepsilon_0 < 1$. Then, for any $f \in L^2(\mathbb{R}^3, C^1(S^2))$, we have

$$\|\mathcal{L}_\varepsilon^c - \sigma_2^2 \Delta_{S^2}\|_{L^2(\mathbb{R}^3 \times S^2)} \leq \varepsilon^{2-\alpha/2} E(f).$$
where, with \( \tilde{f}(x,v) := f(x,v/|v|) \) for \( v \in \mathbb{R}^3 \),

\[
E(f) := \frac{\pi}{3(1 - \varepsilon_0)^6} \sup_{|h| \leq r_0} \lambda \sup_{|h|} \sup_{\alpha} \|D_h^k \tilde{f}(\cdot, + h)\|_{L^2(\mathbb{R}^3 \times \mathbb{S}^2)} + \left( \frac{6\pi}{(1 - \varepsilon_0)^3} \sup_{|h|} \|a''(v)|\right) \sup_{\alpha} \lambda \|\Delta_{\mathbb{S}^2} f\|_{L^2(\mathbb{R}^3 \times \mathbb{S}^2)}.
\]

(34)

**Proof** Before starting the proof, we introduce the retraction \( R_k \) at \( \hat{k} \) onto the sphere \( R_k(h) := \frac{k+h}{|k+h|} \), and

\[
B_{\varepsilon,k} := R_k^{-1}(S_\varepsilon^c) = \{ h = \beta_1 \hat{k}_1 + \beta_2 \hat{k}_2 : \beta = (\beta_1, \beta_2) \in \mathbb{R}^2 \text{ with } |\beta| < r_\varepsilon \},
\]

where \( (\hat{k}_1, \hat{k}_2) \) stands for an orthonormal basis of \( \hat{k}_1 \). We also recall that \( r_\varepsilon = \sqrt{1 - (1 - \varepsilon)^2} \), coming from the relation \( \tan(\arccos(s)) = \sqrt{1 - s^2} / s \) and (34). In different terms, \( B_{\varepsilon,k} \) is a ball centered at \( 0 \) with radius \( r_\varepsilon \) on the tangent plane to the unit sphere at \( \hat{k} \), and the retraction \( R_k \) holds from \( B_{\varepsilon,k} \) onto \( S_\varepsilon^c \).

To prove the lemma, we start with the following change of variables \( \tilde{R} = R_k(h) \) in \( L^2_\varepsilon \), so that

\[
L^2_\varepsilon f(x,\hat{k}) = \frac{\lambda(x)}{2^{1+\alpha(x)/2}} \int_{B_{\varepsilon,k}} \rho(x, R_k(h) \cdot \hat{k}) (f(x, R_k(h)) - f(x, R_k(0))) |\det Jac R_k(h)| dh.
\]

Using that \( \tilde{f}(\hat{x} + h) = f(x, R_k(h)) \) and \( \tilde{f}(x, \hat{k}) = f(x, \hat{k}) \), one can decompose \( L^2_\varepsilon f \) as

\[
L^2_\varepsilon f(x,\hat{k}) = D_1 + D_2 + D_3 + D_4,
\]

where the terms \( D_j \) follow with obvious notations from the Taylor expansion

\[
\tilde{f}(x, \hat{k} + h) - \tilde{f}(x, \hat{k}) = D_k \tilde{f}(x, \hat{k})(h) + \frac{1}{2!} D_k^2 \tilde{f}(x, \hat{k})(h,h) + \frac{1}{3!} D_k^3 \tilde{f}(x, \hat{k})(h,h,h) + \frac{1}{3!} \int_0^1 (1 - s)^3 D_k^3 \tilde{f}(x, \hat{k} + sh)(h,h,h,h) ds.
\]

The terms \( D_1 \) and \( D_3 \). Using that the ball \( B_{\varepsilon,k} \) in the tangent plane is symmetric with respect to \( 0 \), we just make the change of variables \( h \to -h \), so that \( D_1 = -D_1 \) and \( D_3 = -D_3 \) leading to \( D_1 = D_3 = 0 \).

The term \( D_4 \). We have

\[
|D_4| \leq \frac{\lambda(x)}{3!^{1+\alpha(x)/2}} \int_0^1 ds (1 - s) \int_{B_{\varepsilon,k}} dh \rho(x, R_k(h) \cdot \hat{k}) \|D_k^3 \tilde{f}(x, \hat{k} + sh)(h,h,h,h)\| dh.
\]

Since \( R_k(h) = \frac{k+h}{\sqrt{1 + |h|^2}} \), we have

\[
R_k(h) \cdot \hat{k} = \frac{1}{\sqrt{1 + |h|^2}} \quad \text{and} \quad |\det Jac R_k(h)| = \frac{1}{\sqrt{1 + |h|^2}}.
\]

As a consequence, we find

\[
\|D_4\|_{L^2(\mathbb{R}^3 \times \mathbb{S}^2)} \leq \frac{\lambda(x)}{3!^{1+\alpha(x)/2}} \int_{|h| \leq r_\varepsilon} \frac{|h|^4}{(1 - \sqrt{1 + |h|^2})^{1+\alpha(x)/2}} dh.
\]

Changing to polar coordinates in the last integral gives

\[
\int_{|h| \leq r_\varepsilon} \frac{|h|^4}{(1 - \sqrt{1 + |h|^2})^{1+\alpha(x)/2}} dh = 2\pi \int_0^{r_\varepsilon} \frac{v^5}{(1 - \sqrt{1 + v^2})^{1+\alpha(x)/2}} dv
\]

\[
= 2\pi \int_0^{\varepsilon(1 - \varepsilon)} v^{1-\alpha(x)/2} (2 + v)^2 (v + 1)^{2+\alpha(x)/2} dv
\]

where we used the change of variables \( v = \sqrt{1 + r_\varepsilon^2 - 1} = \varepsilon/(1 - \varepsilon) \). This gives finally

\[
\|D_4\|_{L^2(\mathbb{R}^3 \times \mathbb{S}^2)} \leq \frac{\pi}{3} \sup_{|h| \leq r_\varepsilon} \|D_4 \tilde{f}(\cdot, + h)\|_{L^2(\mathbb{R}^3 \times \mathbb{S}^2)} \sup_x \lambda(x) \frac{\varepsilon^{2+\alpha(x)/2}}{(1 - \varepsilon)^6}.
\]

22
The term $D_2$. For this last term, we have

$$D_2 = \frac{\lambda(x)}{2^{2+\alpha(x)/2}} \int_{B_{r',\hat{h}}} \rho(x, R_k(h) \cdot \hat{k}) D^2 \hat{f}(x, \hat{k})(h, h) |\det Jac R_k(h)| dh,$$

with

$$D^2 \hat{f}(x, \hat{k})(h, h) = h_1^2 \partial_{\hat{k}_1 \hat{k}_1} \hat{f}(x, \hat{k}) + h_2^2 \partial_{\hat{k}_2 \hat{k}_2} \hat{f}(x, \hat{k}) + 2h_1 h_2 \partial_{\hat{k}_1 \hat{k}_2} \hat{f}(x, \hat{k}),$$

and, accordingly, the following decomposition $D_2 = D_{21} + D_{22} + 2D_{23}$. Applying the change of variables $h = (h_1, h_2) \rightarrow (-h_1, h_2)$ leads to $D_{23} = 0$. Setting $h = (h_1, h_2) \rightarrow (h_2, h_1)$ leads to

$$D_2 = \frac{\lambda(x)}{2^{2+\alpha(x)/2}} \int_{B_{r',\hat{h}}} \rho(x, R_k(h) \cdot \hat{k}) |h|^2 |\det Jac R_k(h)| dh \text{Trace}(Hess \hat{f}(x, \hat{k})),$$

where

$$\text{Trace}(Hess \hat{f}(x, \hat{k})) = \Delta_p f(x, \frac{p}{|p|})|_{\hat{p}=\hat{k}} = \Delta_{\hat{S}^2} f(x, \hat{k}).$$

Furthermore, with the change of variables $\hat{p} = R_k(h)$, we find

$$\int_{B_{r',\hat{h}}} \rho(x, R_k(h) \cdot \hat{k}) |h|^2 |\det Jac R_k(h)| dh = \int_{S^\epsilon_0^x} \rho(x, \hat{p} \cdot \hat{k}) |R_k^{-1}(\hat{p})|^2 \sigma(d\hat{p}),$$

and note that for $\hat{p} \cdot \hat{k} = \cos(\theta)$, we have $|R_k^{-1}(\hat{p})|^2 = \tan^2(\theta)$. As a result, moving to spherical coordinates, and performing the change of variables $v = \tan(\theta/2)$ together with the relation

$$\arccos(s) = 2 \arctan\left(\frac{\sqrt{1-s^2}}{1+s}\right) \quad \text{for} \quad s \in (-1,1],$$

we find, with $r'_\epsilon = \sqrt{1-(1-\epsilon)^2}/(2-\epsilon)$,

$$\int_{S^\epsilon_0^x} \rho(x, \hat{p} \cdot \hat{k}) |R_k^{-1}(\hat{p})|^2 \sigma(d\hat{p}) = 2\pi \int_0^{\arccos(1-\epsilon)} \rho(x, \cos(\theta)) \tan^2(\theta) \sin(\theta) d\theta$$

$$= 2^{1-\alpha(x)/2} \pi \int_0^{r'_\epsilon} a(2v/\sqrt{1+v^2})(1+v^2)^{\alpha(x)/2} \frac{\sigma(v)}{(1-v^2)^{2\alpha(x)-1}} dv,$$

leading to

$$D_2 = \tilde{D}_2 \Delta_{\hat{S}^2} f(\hat{k}) := 2^{1-\alpha(x)/2} \pi \lambda(x) \int_0^{r'_\epsilon} a(2v/\sqrt{1+v^2})(1+v^2)^{\alpha(x)/2} \frac{\sigma(v)}{(1-v^2)^{2\alpha(x)-1}} dv \Delta_{\hat{S}^2} f(x, \hat{k}).$$

Now, let us introduce

$$\tilde{\sigma}_x^2(\hat{R}) := 2^{1-\alpha(x)/2} \pi \lambda(x) \int_0^{r'_\epsilon} a(2v/\sqrt{1+v^2}) \frac{dv}{v^{\alpha(x)-1}},$$

and remark that

$$|\tilde{D}_2 - \tilde{\sigma}_x^2(\hat{R})| \leq 3 \cdot 2^{3-\alpha(x)/2} \pi \lambda(x) \sup_v a(v) \frac{1}{(1-r'_\epsilon^2)^3} \int_0^{r'_\epsilon} v^{3-\alpha(x)} dv \leq 3 \cdot 2^{2-\alpha(x)/2} \pi \lambda(x) \sup_v a(v) \frac{r'_\epsilon^2}{(1-r'_\epsilon^2)^3} \leq 6\pi \sup_v \lambda(x) \sup_v a(v) \frac{2^{2-\alpha_M/2}}{(1-\epsilon)^3},$$

since $r'_\epsilon \leq \sqrt{2\epsilon}$, $0 \leq \alpha_m \leq \alpha(x) \leq \alpha_M < 2$ and $1-r'_\epsilon^2 = \frac{2(1-\epsilon)}{2-\epsilon} > 2(1-\epsilon)$. With the definition of $\sigma_\epsilon$ given in (10), we obtain using $a'(0) = 0$,

$$|\tilde{\sigma}_x^2(\hat{R}) - \sigma_\epsilon^2(\hat{R})| \leq 2^2 \pi \lambda(x) \sup_v a(v) \left|\int_0^{r'_\epsilon} v^{3-\alpha(x)} dv, \right.$$

$$\leq 2^2 \pi \sup_v \lambda(x) \sup_v a(v) \frac{2^{2-\alpha_M/2}}{(1-\epsilon)^3}. $$

Collecting the various estimates on the $D_j$ and using that $\epsilon \leq \epsilon_0 < 1$ concludes the proof of Lemma 6.1 and therefore of Proposition 2.1 □
6.2 Proof of Proposition 2.2

We first show that the infinitesimal generator of the Markov process $D_z$ is $A_z$.

**Infinitesimal generator for $D_z$.** Let $f$ be a smooth bounded function on $\mathbb{R}^3 \times \mathbb{S}^2$. The goal of this section is to prove that

$$\lim_{h \to 0^+} \frac{1}{h} (\mathbb{E}_z[f(D_z(h))] - f(z)) = A_z f(z). \tag{35}$$

To this end, we introduce the first jump time $T_1$ to obtain

$$\mathbb{E}_z[f(D_z(t))] = \mathbb{E}_z[f(D_z(t))1_{(T_1 > t)}] + \mathbb{E}_z[f(D_z(t))1_{(T_1 \leq t)}]. \tag{36}$$

Using conditional expectations, we find for the first term

$$\mathbb{E}_z[f(D_z(t))1_{(T_1 > t)}] = \mathbb{E}_z[f(\psi_0^z(t))1_{(T_1 > t)}] = \mathbb{E}_z[\mathbb{E}_z[f(\psi_0^z(t))1_{(T_1 > t)} | \psi_0^z(s), s \in [0, t])]
\begin{align*}
&= \mathbb{E}_z\left[f(\psi_0^z(t)) \mathbb{P}_z(T_1 > t | \psi_0^z(s), s \in [0, t])\right]
&= \mathbb{E}_z\left[f(\psi_0^z(t)) e^{-\int_0^t A_z(\psi_0^z(s))ds}\right].
\end{align*}

With the following notations for the flow $\psi_n^z$,

$$\psi_n^z = (X_n^z, K_n^z) = ((X_{j,n}^z)_{j=1,2,3}, (K_{j,n}^z)_{j=1,2,3}) \in \mathbb{R}^3 \times \mathbb{S}^2,$$

together with $f(x, k) = f(x, k/|k|)$ for $(x, k) \in \mathbb{R}^3 \times \mathbb{R}^3$, the Itô formula yields

$$df(\psi_n^z(t)) = d\tilde{f}(\psi_n^z(t)) = \nabla_x \tilde{f}(\psi_n^z(t)) \cdot dX_n^z(t) + \nabla_k \tilde{f}(\psi_n^z(t)) \cdot dK_n^z(t) + \frac{1}{2} \sum_{j,k=1,2,3} \tilde{K}_{j,k} \tilde{f}(\psi_n^z(t)) d < K_{j,n}^z, K_{k,n}^z > + |\tilde{f}(\psi_n^z(t))| \Delta_{S^2} f(\psi_n^z(t)) dt$$

Above, we have used the fact that

$$\Delta_{S^2} f(x, \hat{k}) = \Delta_k \tilde{f}(x, \hat{k}) - \sum_{j,k=1,2,3} \hat{k}_j \hat{k}_k \partial_{k_j}^2 \tilde{f}(x, \hat{k}) - 2 \sum_{j,k=1,2,3} \hat{k}_j \partial_{k_j} \tilde{f}(x, \hat{k}). \tag{37}$$

Therefore, we have for $n = 0$,

$$d\left((f(\psi_0^z(t)) - f(z)) e^{-\int_0^t A_z(\psi_0^z(s))ds}\right) = \left(df(\psi_0^z(t)) - A_z(\psi_0^z(t))(f(\psi_0^z(t)) - f(z))\right) e^{-\int_0^t A_z(\psi_0^z(s))ds}$$

so that

$$\lim_{h \to 0^+} \frac{1}{h} \mathbb{E}_z[(f(D_z(h)) - f(z))1_{(T_1 > h)}] = \hat{k} \cdot \nabla_x f(z) + \sigma_x^2(\epsilon) \Delta_{S^2} f(z), \quad \forall z = (x, \hat{k}) \in \mathbb{R}^3 \times \mathbb{S}^2.$$

Regarding the second term in (36), we find, using the Markov property in the third line,

$$\mathbb{E}_z[f(D_z(h))1_{(T_1 \leq h)}] = \mathbb{E}_z[\mathbb{E}_z[f(D_z(h)) | T_1] 1_{(T_1 \leq h)}]$$

$$= \mathbb{E}_z\left[\int_0^{h} \mathbb{E}_z[f(D_z(h)| T_1 = v)] \Lambda_z(\psi_0^z(v)) e^{-\int_0^v A_z(\psi_0^z(s))ds} dv \right]$$

$$= \mathbb{E}_z\left[\int_0^{h} \int_{\mathbb{R}^3 \times \mathbb{S}^2} \mathbb{E}_{\nu}[f(D_z(h-v))] \Lambda_z(\psi_0^z(v), d\nu) \Lambda_z(\psi_0^z(v)) e^{-\int_v^h A_z(\psi_0^z(s))ds} dv \right]$$

$$= \mathbb{E}_z\left[\int_0^{h} \int_{\mathbb{S}^2} \mathbb{E}[\chi_{\hat{z} = (h-w, \hat{p})} f(D_z(w))] \pi_z(\psi_0^z(h-w), \hat{p}) \sigma(\hat{p}) \Lambda_z(\psi_0^z(h-w)) e^{-\int_0^{h-w} A_z(\psi_0^z(s))ds} dw \right],$$

where the probability $\Pi_z$ and the density $\pi_z$ are defined respectively in (16) and (17). As a consequence,

$$\lim_{h \to 0^+} \frac{1}{h} \mathbb{E}_z[f(D_z(h))1_{(T_1 \leq h)}] = \Lambda_z(z) \int_{\mathbb{S}^2} f(x, \hat{p}) \pi_z(z, \hat{p}) \sigma(\hat{p}).$$
Moreover, we have
\[ \mathbb{P}_z(T_1 \leq h) = \mathbb{E}_z \left[ \mathbb{P}_z(T_1 \leq h \mid \psi_0^z(s), s \in [0, h]) \right] = \mathbb{E}_z \left[ 1 - e^{-\int_0^h \lambda_s(\psi_0^z(s))ds} \right], \]
and it is then direct to see that
\[ \lim_{h \to 0} \frac{1}{h} \mathbb{P}_z(T_1 \leq h) = \Lambda_z(z). \]
This finally yields
\[ \lim_{h \to 0^+} \frac{1}{h} \mathbb{E}_z[(f(D_z(h)) - f(z))1_{(T_1 \leq h)}] = \Lambda_z(z) \int_{\mathbb{S}^2} (f(x, \tilde{\nu}) - f(x, \tilde{\nu})) \pi_z(z, \tilde{\nu}) \sigma(d\tilde{\nu}), \]
which gives (35) collecting all results.

**Proof of (38).** Since \( D_z \) is a solution to the martingale problem associated to \( \mathcal{A}_z \) (see [15, Proposition 1.7, pp. 162]), we have, for any smooth bounded function \( f \) on \( \mathbb{R}^3 \times \mathbb{S}^2 \),
\[ \mu_t(f) = \mu_0(f) + \int_0^t \mu_s(\mathcal{A}_z f) ds, \quad \text{where} \quad \mu_t(f) := \mathbb{E}_{\mu_0}[f(D_z(t))]. \]
Let
\[ \nu_t(f) := \frac{1}{u_0} \int_{\mathbb{R}^3 \times \mathbb{S}^2} u_z(t, x, \tilde{k}) f(x, \tilde{k}) dx \sigma(d\tilde{k}). \]
Since \( u_z \) solves (9), we have
\[ u_z(t) = u_0 + \int_0^t \mathcal{A}_z^* u_z(s) ds, \]
where \( \mathcal{A}_z^* \) stands for the adjoint operator of \( \mathcal{A}_z \) in \( L^2(\mathbb{R}^3 \times \mathbb{S}^2) \). Then,
\begin{align*}
\nu_t(f) &= \mu_0(f) + \frac{1}{u_0} \int_{\mathbb{R}^3 \times \mathbb{S}^2} \mathcal{A}_z^* u_z(s, x, \tilde{k}) ds f(x, \tilde{k}) dx \sigma(d\tilde{k}) \\
&= \mu_0(f) + \int_0^t \frac{1}{u_0} \int_{\mathbb{R}^3 \times \mathbb{S}^2} u_z(s, x, \tilde{k}) \mathcal{A}_z f(x, \tilde{k}) dx \sigma(d\tilde{k}) ds = \mu_0(f) + \int_0^t \nu_s(\mathcal{A}_z f) ds.
\end{align*}
Therefore, according to [15, Proposition 9.18, pp. 251], we have \( \mu_t = \nu_t \) for any \( t \geq 0 \), which concludes the proof.

### 6.3 Proof of Theorem 3.1

The proof of this result is provided in three steps. The first step consists in rewriting the probabilistic representation (38) for (9) in terms of a SDE with jumps. The second step concerns the error analysis of the solution to this later SDE with its discretized version. Finally, the last step gathers all the error estimated and concludes the proof.

#### 6.3.1 Step 1

We first introduce an equivalent formulation (in the statistical sense) for the process \( D_z \) in terms of a stochastic differential equation (SDE) with jumps. This representation is useful when comparing with the discrete scheme. Let \( \tilde{D}_z = (\tilde{X}_z, \tilde{K}_z) \) be the solution to the following SDE with jumps:
\begin{align*}
d\tilde{X}_z(t) &= \tilde{K}_z(t) dt \\
d\tilde{K}_z(t) &= \sqrt{2} \sigma_z(\tilde{X}_z(t^-)) \tilde{K}_z(t^-) \times dW - 2 \sigma_z^2(\tilde{X}_z(t^-)) \tilde{K}_z(t^-) dt \\
&\quad + \int_{(0, \pi) \times (0, 2\pi) \times (0, 1)} \tilde{R}(\theta, \varphi, \tilde{D}_z(t^-), v) P(\theta, \theta, \varphi, \varphi, v) (dt, d\theta, d\varphi, dv)
\end{align*}
(38)
where the function \( \tilde{R} \) is defined by, for \( z = (x, \tilde{k}) \),
\[ \tilde{R}(\theta, \varphi, z, v) = \begin{cases} 
R(\theta, \varphi, \tilde{k}) - \tilde{k} & \text{if } v \leq \Lambda_z(z) \pi_z(z, R(\theta, \varphi, \tilde{k})) / \Lambda_z \\
0 & \text{otherwise}.
\end{cases} \]
Above, $R$ is defined in (28), $A_{\varepsilon}$ in (14), $\tilde{A}_{\varepsilon}$ in (21), $\pi_{\varepsilon}$ in (17), and $P$ is a random Poisson measure with intensity measure

$$\mu(dt,d\theta,d\varphi,du) = \Lambda_{\varepsilon}1_{(0,\infty)\times(0,\pi)\times(0,2\pi)\times(0,1)}(t,\varphi,v)\sin(\theta)dtd\varphi d\theta/(4\pi).$$

(39)

See e.g. [2] Chapter 2 for more details on Poisson random measures. The notation $t^-$ is standard and refers to the left limit when approaching $t$ before a jump. With this construction, the infinitesimal generator for the Markov process $\tilde{D}_{\varepsilon}$ is

$$\mathcal{A}_{\varepsilon} f(z) = \hat{k} \cdot \nabla_x f(z) + \sigma_x^2(x)\Delta_{\varepsilon} f(z) + \tilde{A}_{\varepsilon} \int_{(0,\pi)\times(0,2\pi)\times(0,1)} (f(x, \hat{k} + \hat{R}(\theta, \varphi, z, v)) - f(z))\sigma(d\hat{p})dv,$$

and we have the following result.

**Lemma 6.2** The processes $\tilde{D}_{\varepsilon}$ and $D_{\varepsilon}$ have the same distribution.

**Proof** $\tilde{D}_{\varepsilon}$ and $D_{\varepsilon}$ are both Markov processes and are therefore characterized by their generators. We just have then to prove that $\mathcal{A}_{\varepsilon} g = \tilde{A}_{\varepsilon} g$ for any bounded smooth function $g$. This is a direct consequence of the definition to $\hat{R}$. Indeed, denoting by $\hat{I}_{\varepsilon}$ and $\tilde{I}_{\varepsilon}$ the integral operators in respectively $\mathcal{A}_{\varepsilon}$ and $\tilde{A}_{\varepsilon}$, we have with $z = (x, k)$,

$$\tilde{I}_{\varepsilon} g(z) = \tilde{A}_{\varepsilon} \int_{(0,\pi)\times(0,2\pi)\times(0,1)} (g(x, \hat{k} + \hat{R}(\theta, \varphi, z, v)) - g(z))\sin(\theta) d\varphi d\theta dv/(4\pi)$$

$$= \tilde{A}_{\varepsilon} \int_{(0,\pi)\times(0,2\pi)\times(0,1)} \int_0^{\Lambda_{\varepsilon}(z)\pi_{\varepsilon}(z, R(\theta, \varphi, k))} d\hat{p} (g(x, R(\theta, \varphi, \hat{k})) - g(z))\sin(\theta) d\varphi d\theta dv/(4\pi)$$

$$= \tilde{A}_{\varepsilon} \int_{\mathbb{R}^3} (f(x, \hat{k}) - g(z))\sigma(d\hat{p}) = \tilde{I}_{\varepsilon} g(z),$$

which concludes the proof. $\square$

### 6.3.2 Step 2

The goal is now to prove that the discretized process $D_{h,\varepsilon}$ approximates $\tilde{D}_{\varepsilon}$ in a statistical sense. We use for this the notations of Section 3.2 for $X_{n,m}$, $K_{n,m}$, and $K_{n,m}$. For simplicity, we suppose that the Gaussian vectors $(W_{n,m})$ in (22) are obtained from a single 3D standard Brownian motion $W$ as follows: for $n \geq 0$ and $m = 0, \ldots, m_n$, we set $W_{n,m} = (W(T_{n,m} + h_{n,m}) - W(T_{n,m}))/\sqrt{h_{n,m}}$. In the sequel, we will use the following process, defined by, for $t \in [T_{n,m}, T_{n,m} + h_{n,m}]$,

$$\begin{cases}
X_{n,m}(t) = X_{n,m} + \int_{T_{n,m}}^t \tilde{K}_{n,m}ds \\
K_{n,m}(t) = \tilde{K}_{n,m} - 2\int_{T_{n,m}}^t \sigma_{\varepsilon}(X_{n,m})\tilde{K}_{n,m}ds + \sqrt{2} \int_{T_{n,m}}^t \sigma_{\varepsilon}(X_{n,m})\tilde{K}_{n,m}dW(s).
\end{cases}
$$

(40)

For $t \geq 0$, we then combine the $(X_{n,m}, K_{n,m})$ into

$$\mathcal{D}_{h,\varepsilon}(t) = (X_{h,\varepsilon}(t), K_{h,\varepsilon}(t)) = \sum_{n=0}^\infty \sum_{m=0}^{m_n} 1_{[T_{n,m}, T_{n,m} + h_{n,m}]}(t)\Psi_{n,m}(t),$$

where $\Psi_{n,m}(t) = (X_{n,m}(t), K_{n,m}(t))$ is the solution to (40) with initial condition $\Psi_{n,m}(T_{n,m}) = (X_{n,m}, \tilde{K}_{n,m})$. Note that $\mathcal{D}_{h,\varepsilon}$ is simply an interpolation of $D_{h,\varepsilon}$ in the intervals $[T_{n,m}, T_{n,m} + h_{n,m}]$ that will allow us to use the Itô formula.

For any smooth function $f$, we now introduce $b$ (the smooth) solution to the following backward RTE,

$$\partial_t b + \hat{k} \cdot \nabla_x b + \frac{\lambda(x)}{2\sqrt{2\pi \alpha(x)}} b - \rho(x, \hat{k})\sigma_d(d\hat{p}) = 0,$$

(41)

with terminal condition $b(T, x, k) = f(x, \hat{k})$ and use the notation $\hat{b}(t, x, k) = b(t, x, k/|k|)$, $(t, x, k) \in [0, T] \times \mathbb{R}^3 \times \mathbb{R}^3$. We have the following result.

**Proposition 6.1** For any $T > 0$, any smooth bounded function $f$ on $\mathbb{R}^3 \times \mathbb{S}^2$, and any $(x, \hat{k}) \in \mathbb{R}^3 \times \mathbb{S}^2$, we have

$$|E_{(x,\hat{k})}[f(D_{h,\varepsilon}(T))] - E_{(x,\hat{k})}[f(\tilde{D}_{\varepsilon}(T))]| \leq \varepsilon^2 \alpha T^{1/2} 2T E_{\infty}(b) + hF_1(b),$$

where $E_{\infty}(b)$ is defined as in (34) with $L^2$ norms replaced by $L^\infty$ norms in all variables, and where $F_1(b)$ is an explicit function independent of $\varepsilon$ and $h$ that depends on derivatives of $b$ up to order 4.
The notation $\mathbb{E}_{(x,k)}$ above indicates that the process under the expectation starts at the point $(x, \hat{k})$.

**Proof** The proof consists in analyzing the discretization error of the diffusion process in a weak sense following the ideas of [53].

We first write

$$E_{h,\varepsilon} := \mathbb{E}_z[f(D_{h,\varepsilon}(T))] - \mathbb{E}_z[f(\hat{D}_z(T))] = \mathbb{E}_z[b(T, D_{h,\varepsilon}(T))] - \mathbb{E}_z[b(T, \hat{D}_z(T))],$$

with $z = (x, \hat{k})$. Using the Itô formula, see e.g. [2] Th 4.4.7 p. 226, together with the expression of $\Delta_{S\varepsilon}$ given in [37], we have

$$\mathbb{E}_z[b(T, D_z(T))] - b(0, z) = \mathbb{E}_z\left[\int_0^T \partial_t b(t, \hat{D}_z(t)) + \hat{K}_z(t) \cdot \nabla_x \varepsilon b(t, \hat{D}_z(t)) + \sigma^2(\hat{X}_z(t)) \Delta_{S\varepsilon} b(t, \hat{D}_z(t)) + T^\varepsilon(t) dt\right],$$

where

$$T^\varepsilon(t) = \frac{\lambda(\hat{X}_z(t))}{21 + \alpha(\hat{X}_z(t))} \int_{S_{\varepsilon}} \rho(\hat{X}_z(t), \hat{K}_z(t) \cdot \hat{p})(b(t, \hat{X}_z(t), \hat{p}) - b(t, \hat{X}_z(t), \hat{K}_z(t))\sigma(d\hat{p}).$$

Using the fact that $b$ satisfies [41], we find

$$\mathbb{E}_z[b(T, D_z(T))] = b(0, z) + \mathbb{E}_z\left[\int_0^T \sigma^2(\hat{X}_z(t)) \Delta_{S\varepsilon} b(t, \hat{D}_z(t)) - T^\varepsilon(t) dt\right],$$

where $T^\varepsilon$ is as $T^\varepsilon$ with $S_{\varepsilon}$ replaced by $S^\varepsilon$. Following the lines of the proof of Lemma 6.1, we obtain

$$|\mathbb{E}[b(T, D_z(T))] - b(0, z)| \leq \varepsilon^{2-\alpha} M/2 T E_\infty(b),$$

where $E_\infty(b)$ is defined as in [54] with $L^2$ norms replaced by $L^\infty$ norms. We move now to the term $\mathbb{E}_z[b(T, D_{h,\varepsilon}(T))]$ which requires more work. Decomposing the interval $[0, T]$ according to the grid $(T_{n,m})$, we have

$$b(T, D_{h,\varepsilon}(T)) - b(0, z) = \hat{b}(T, D_{h,\varepsilon}(T)) - \hat{b}(0, z)$$

$$= \sum_{n \geq 0} \sum_{m \geq 0} \hat{b}(T_{n,m+1} \wedge T, D_{h,\varepsilon}(T_{n,m+1} \wedge T)) - \hat{b}(T_{n,m} \wedge T, D_{h,\varepsilon}(T_{n,m} \wedge T))$$

$$+ \hat{b}(T_{n,m} \wedge T, D_{h,\varepsilon}(T_{n,m} \wedge T)) - \hat{b}(T_{n,m} \wedge T, D_{h,\varepsilon}(T_{n,m} \wedge T^-))$$

$$=: B_1 + B_2,$$

with obvious notations and where $B_1$ is meant to capture the dynamic between jumps while $B_2$ captures that at the jumps. The double sum and the $T_{n,m} \wedge T$ are only here to simplify the proof. Note that in order to define the sum for all $m \geq 0$, we set $T_{n,m} = T$ for $m > m_n$, and note also that there is only a finite number of terms in the sums. We are then led to estimate the differences in $B_1$ and $B_2$ for which we will use the process $D_{h,\varepsilon}$. We introduced $\hat{b}$ since $\hat{r}_{h,\varepsilon}$ is not necessarily on the sphere between the grid points. Since $T$ is on the grid, we have by definition $D_{h,\varepsilon}(T_{n,m} \wedge T) = D_{h,\varepsilon}(T_{n,m} \wedge T)$. Consider now the notation

$$\hat{D}_{h,\varepsilon} := (\hat{X}_{h,\varepsilon}, \hat{R}_{h,\varepsilon}) \quad \text{with} \quad \hat{R}_{h,\varepsilon} := \frac{\hat{r}_{h,\varepsilon}}{\hat{r}_{h,\varepsilon}},$$

By construction, the process $\hat{D}_{h,\varepsilon}$ is continuous at the times $T_{n,m}$ that do not correspond to jump times, so that

$$\hat{b}(T_{n,m}, \hat{D}_{h,\varepsilon}(T_{n,m})) = b(T_{n,m}, \hat{D}_{h,\varepsilon}(T_{n,m})) = b(T_{n,m}, \hat{D}_{h,\varepsilon}(T_{n,m}^-)) = \hat{b}(T_{n,m}, \hat{D}_{h,\varepsilon}(T_{n,m}^-))$$

for those $T_{n,m}$. As a consequence, $B_2$ indeed only accounts for jumps. We will then estimate $B_1$ using the Itô formula for $D_{h,\varepsilon}$ between $T_{n,m}$ and $T_{n,m+1}$, and the properties of Poisson random measures for $B_2$. For the latter, we notice that we have by construction $D_{h,\varepsilon}(T_n) = \hat{D}_{h,\varepsilon}(T_n^-)$. Using then the random Poisson measure $P$ with intensity measure $\mu$ introduced in [38] and [59], we can write

$$B_2 = \int_0^T ((\hat{b}(t, \hat{D}_{h,\varepsilon}(t^-)) + \hat{R}(\theta, \varphi, \hat{D}_{h,\varepsilon}(t^-), v) - \hat{b}(t, \hat{D}_{h,\varepsilon}(t^-))) P(dt, d\theta, d\varphi, dv),$$

with obvious notations.
so that, together with the fact that $P - \mu$ is a measure-valued martingale, see e.g. [2 Chapter 2],

$$
\mathbb{E}_z[B_2] = \sum_n \sum_m \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \tilde{I}^z_>(t) dt \right]
$$

with

$$
\tilde{I}^z_>(t) = \frac{\lambda(X_{b,x}(t))}{2^{1+\alpha(X_{b,x}(t))/2}} \int_{S^2} \rho(X_{b,x}(t), \hat{R}_{b,x}(t) \cdot \hat{\rho})(\hat{b}(t, X_{b,x}(t), \hat{\rho}) - \hat{b}(t, X_{b,x}(t), \hat{R}_{b,x}(t))\sigma(d\hat{\rho}).
$$

For $B_1$, we have from the Itô formula

$$
\mathbb{E}_z[B_1] = \sum_n \sum_m \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \partial_t \hat{b}(t, \mathcal{D}_{h,e}(t)) + \hat{R}_{h,e}(T_{n,m}) \cdot \nabla \hat{b}(t, \mathcal{D}_{h,e}(t)) + \sigma_0^2(X_{b,x}(T_{n,m})) \mathcal{B}^e_{n,m}(t) dt \right]
$$

where

$$
\mathcal{B}^e_{n,m}(t) = \Delta_k \hat{b}(t, \mathcal{D}_{h,e}(t)) - \mathcal{R}_{h,e}(T_{n,m})^T D^2 \hat{b}(t, \mathcal{D}_{h,e}(t)) \mathcal{R}_{h,e}(T_{n,m}) - 2 \mathcal{R}_{h,e}(T_{n,m}) \cdot \nabla_k \hat{b}(t, \mathcal{D}_{h,e}(t)).
$$

As a result,

$$
\mathbb{E}_z[B_1 + B_2] = \sum_n \sum_m \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \partial_t \hat{b}(t, \mathcal{D}_{h,e}(t)) + \hat{R}_{h,e}(T_{n,m}) \cdot \nabla \hat{b}(t, \mathcal{D}_{h,e}(t)) + \sigma_0^2(X_{b,x}(T_{n,m})) \mathcal{B}^e_{n,m}(t) + \tilde{I}^z_>(\mathcal{D}_{h,e}(t)) dt \right]
$$

Using again the fact that $b$ satisfies (41), we have

$$
\partial_t \hat{b}(t, \mathcal{D}_{h,e}(T_{n,m})) + \hat{R}_{h,e}(T_{n,m}) \cdot \nabla \hat{b}(t, \mathcal{D}_{h,e}(T_{n,m})) + \mathcal{Q}b(t, \mathcal{D}_{h,e}(T_{n,m})) = 0,
$$

which also holds true for $\hat{b}$ since the variable $\mathcal{R}_{h,e}(T_{n,m})$ has norm 1 at the grid points $T_{n,m}$. As a consequence,

$$
\mathbb{E}_z[b(T, D_{h,e}(T))] - b(0, z) = \sum_n \sum_m \sum_{j=1}^7 E_{n,m}^j
$$

where

$$
E_{n,m}^{(1)} := \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \partial_t \hat{b}(t, \mathcal{D}_{h,e}(t)) - \partial_t \hat{b}(t, \mathcal{D}_{h,e}(T_{n,m})) dt \right]
$$

$$
E_{n,m}^{(2)} := \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \hat{R}_{h,e}(T_{n,m}) \cdot (\nabla_k \hat{b}(t, \mathcal{D}_{h,e}(t)) - \nabla_k \hat{b}(t, \mathcal{D}_{h,e}(T_{n,m}))) dt \right]
$$

$$
E_{n,m}^{(3)} := \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \sigma_0^2(X_{b,x}(T_{n,m})) (\Delta_k \hat{b}(t, \mathcal{D}_{h,e}(t)) - \Delta_k \hat{b}(t, \mathcal{D}_{h,e}(T_{n,m}))) dt \right]
$$

$$
E_{n,m}^{(4)} := \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} -\sigma_0^2(X_{b,x}(T_{n,m})) \mathcal{R}_{h,e}(T_{n,m})^T D^2 \hat{b}(t, \mathcal{D}_{h,e}(t)) - D^2 \hat{b}(t, \mathcal{D}_{h,e}(T_{n,m}))) \mathcal{R}_{h,e}(T_{n,m}) dt \right]
$$

$$
E_{n,m}^{(5)} := \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} -2 \sigma_0^2(X_{b,x}(T_{n,m})) \mathcal{R}_{h,e}(T_{n,m}) \cdot (\nabla_k \hat{b}(t, \mathcal{D}_{h,e}(t)) - \nabla_k \hat{b}(t, \mathcal{D}_{h,e}(T_{n,m}))) dt \right]
$$

$$
E_{n,m}^{(6)} := \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \sigma_0^2(X_{b,x}(T_{n,m})) \Delta_k \mathcal{B}(t, \mathcal{D}_{h,e}(T_{n,m})) - \tilde{I}^z_>(\mathcal{D}_{h,e}(T_{n,m})) dt \right]
$$

$$
E_{n,m}^{(7)} := \mathbb{E}_z \left[ \int_{T_{n,m} \wedge T}^{T_{n,m+1} \wedge T} \tilde{I}^z_>(\mathcal{D}_{h,e}(t)) - \tilde{I}^z_>(\mathcal{D}_{h,e}(T_{n,m})) dt \right]
$$

with

$$
\tilde{I}^z_>(\mathcal{D}_{h,e}(t)) = \frac{\lambda(X_{b,x}(t))}{2^{1+\alpha(X_{b,x}(t))}} \int_{S^2} \rho(X_{b,x}(t), \hat{R}_{b,x}(t) \cdot \hat{\rho})(\hat{b}(t, X_{b,x}(t), \hat{\rho}) - \hat{b}(t, X_{b,x}(t), \hat{R}_{b,x}(t))\sigma(d\hat{\rho}).
$$

In the estimates below involving $\hat{b}$, derivatives involving $k \mapsto \hat{k} = k/|k|$ produce terms of the form $1/|k|^p$ for some $p > 0$. These terms are due to the fact that $\mathcal{R}_{n,m}$ does not stay on the sphere at all times. However, these terms can be bounded uniformly thanks to the following lemma.
Lemma 6.3 We have for any $n \geq 0$ and $m \leq m_n$,

$$\inf_{s \in [T_{n,m}, T_{n,m} + h_{n,m}]} |K_{n,m}(s)| \geq (1 - 2h\sigma^2_{\varepsilon,\infty}) \geq \frac{1}{2},$$

for $h$ and $\varepsilon$ small enough, and

$$\sup_{s \in [T_{n,m}, T_{n,m} + h_{n,m}]} \mathbb{E}[|K_{n,m}(s)|^2 | T_{n,m}, T_{n,m+1}] \leq 1 + 4h\sigma^2_{\varepsilon,\infty} \leq 2.$$

Proof For $s \in [T_{n,m}, T_{n,m} + h_{n,m}]$, $K_{n,m}$ can be rewritten as the sum of two orthogonal components

$$K_{n,m}(s) = (1 - 2(t - T_{n,m})\sigma_\varepsilon^2(X_{n,m}))\hat{K}_{n,m} + \sqrt{2}\sigma_\varepsilon(X_{n,m})\hat{K}_{n,m} \times (W(s) - W(T_{n,m}))$$

so that

$$|K_{n,m}(s)|^2 \geq (1 - 2h\sigma^2_{\varepsilon,\infty})^2.$$

Now for the upper bound, using that $W_n$ and $P$ are independent, we have

$$\mathbb{E}[|K_{n,m}(s)|^2 | T_{n,m}, T_{n,m+1}] \leq (1 + 2h\sigma^2_{\varepsilon,\infty}) + 2\sigma^2_{\varepsilon,\infty} \mathbb{E}[|W(s) - W(T_{n,m})|^2 | T_{n,m}, T_{n,m+1}]$$

$$\leq 1 + 4h\sigma^2_{\varepsilon,\infty} \leq 2,$$

where we used that $4h\sigma^2_{\varepsilon,\infty} \leq 1$. This concludes the proof. \(\square\)

For $E_{n,m}^{(1)}$. Using the Itô formula between $T_{n,m}$ and $T_{n,m+1}$, we find

$$\mathbb{E}\left[\partial_t b(t, \mathcal{D}_{n,m}(t)) - \partial_t \hat{b}(t, \mathcal{D}_{n,m}(T_{n,m})) \middle| T_{n,m}, T_{n,m+1}\right]$$

$$= \int_{T_{n,m}}^t \mathbb{E}\left[\partial_t \hat{b}(t, \mathcal{D}_{n,m}(s)) + \sigma_\varepsilon^2(X_{n,m}(T_{n,m})) \left(\Delta_k \partial_t \hat{b}(t, \mathcal{D}_{n,m}(s)) - \mathbb{E}[\hat{K}_{n,m}(T_{n,m})] T \mathbb{E}[\partial_t \hat{b}(t, \mathcal{D}_{n,m}(s))]\right) \right] ds,$$

so that

$$|E_{n,m}^{(1)}| \leq h \left(\|D_{n,m}^2 \hat{b}\|_\infty + 2\sigma^2_{\varepsilon,\infty} \left(\|D_{n,m}^3 \hat{b}\|_\infty + \|D_{n,m}^2 \hat{b}\|_\infty\right)\right),$$

where $\sigma_{\infty,\varepsilon} = \sup_x \sigma_\varepsilon(x)$. When $\varepsilon \leq \varepsilon_0$, we have $\sigma_{\infty,\varepsilon} \leq \sigma_{\infty,\varepsilon_0}$ since $r_\varepsilon'$ is an increasing function of $\varepsilon$. As a result, we obtain

$$\sum_n \sum_m |E_{n,m}^{(1)}| \leq h T \left(\|D_{n,m}^2 \hat{b}\|_\infty + 2\sigma^2_{\varepsilon,\infty} \left(\|D_{n,m}^3 \hat{b}\|_\infty + \|D_{n,m}^2 \hat{b}\|_\infty\right)\right) =: hF_2(b).$$

For $E_{n,m}^{(j)}$ with $j = 2, \ldots, 5$. Following the same lines as above, we have

$$\sum_n \sum_m \sum_{j=2}^5 |E_{n,m}^{(j)}| \leq 2h T \left(\|D_{n,m}^2 \hat{b}\|_\infty + 2\sigma^2_{\varepsilon,\infty} \left(\|D_{n,m}^3 \hat{b}\|_\infty + \|D_{n,m}^2 \hat{b}\|_\infty\right)\right)$$

$$+ 4h T \sigma^2_{\varepsilon,\infty} \left(\|D_{n,m}^3 \hat{b}\|_\infty + \|D_{n,m}^2 \hat{b}\|_\infty + 2\sigma^2_{\varepsilon,\infty} \left(\|D_{n,m}^3 \hat{b}\|_\infty + \|D_{n,m}^2 \hat{b}\|_\infty\right)\right)$$

$$=: hF_3(b).$$

For $E_{n,m}^{(6)}$. For this term, we follow the proof of Lemma 6.1 and find

$$\sum_n \sum_m |E_{n,m}^{(6)}| \leq T \varepsilon^{2-\alpha/2} E_\infty(b),$$

where $E_\infty(b)$ is defined as before.
For $E_{n,m}^{(7)}$, starting from (42), we have
\[
\hat{T}_x^\varepsilon (\Omega_{h,s}(t)) = \frac{\lambda(\Omega_{h,s}(t))}{21+\sigma(\Omega_{h,s}(t))\varepsilon/2} \int_{S^2} \rho(\Omega_{h,s}(t), \hat{\Omega}_{h,s}(t)) (\partial \cdot \hat{\rho})(\hat{\Omega}_{h,s}(t)) - \hat{\rho}(\hat{\Omega}_{h,s}(t), \hat{\rho})(\hat{\rho})(\hat{\Omega}_{h,s}(t)) \sigma(d\hat{\rho}) \\
= \frac{\lambda(\Omega_{h,s}(t))}{21+\sigma(\Omega_{h,s}(t))\varepsilon/2} \int_{0}^{2\pi} d\varphi \int_{-1}^{1-\varepsilon} ds \rho(\Omega_{h,s}(t), s) \\
\times \left( \hat{b}(t, \Omega_{h,s}(t), s, \hat{\Omega}_{h,s}(t) + \sqrt{1-s^2} G(\varphi, \hat{\Omega}_{h,s}(t))) - \hat{b}(t, \Omega_{h,s}(t), \hat{\Omega}_{h,s}(t)) \right) 
\]
where the last line is obtained by changing to spherical coordinates with $s = \cos(\theta)$, and
\[G(\varphi, \kappa) := \cos(\varphi)k_1^1 + \sin(\varphi)k_2^1.
\]
Above, $(k_1^1, k_2^1)$ forms an orthonormal basis of the plane $\hat{k}^\perp$. Note that the choice of $(k_1^1, k_2^1)$ does not play any role since the variable $\varphi$ is integrated. Now, writing
\[
\hat{T}_x^\varepsilon (\Omega_{h,s}(t)) = \frac{\lambda(\Omega_{h,s}(t))}{21+\sigma(\Omega_{h,s}(t))\varepsilon/2} \int_{0}^{2\pi} d\varphi \int_{-1}^{1-\varepsilon} ds \int_{0}^{1} dv \rho(\Omega_{h,s}(t), s) \\
\times \left( (s-1)\hat{\Omega}_{h,s}(t) + \sqrt{1-s^2} G(\varphi, \hat{\Omega}_{h,s}(t)) \right) \\
\cdot \nabla_k \hat{b}(t, \Omega_{h,s}(t), (1+v(s-1))\hat{\Omega}_{h,s}(t) + v\sqrt{1-s^2} G(\varphi, \hat{\Omega}_{h,s}(t)))
\]
and using that $G(\varphi + \pi, \kappa) = -G(\varphi, \kappa)$, we just have to focus on
\[
\hat{T}_x^\varepsilon (\Omega_{h,s}(t)) = \frac{\lambda(\Omega_{h,s}(t))}{21+\sigma(\Omega_{h,s}(t))\varepsilon/2} \int_{0}^{2\pi} d\varphi \int_{-1}^{1-\varepsilon} ds \int_{0}^{1} dv \rho(\Omega_{h,s}(t), s(s-1))\hat{\Omega}_{h,s}(t) \\
\cdot \nabla_k \hat{b}(t, \Omega_{h,s}(t), (1+v(s-1))\hat{\Omega}_{h,s}(t) + v\sqrt{1-s^2} G(\varphi, \hat{\Omega}_{h,s}(t)))
\]
Before applying the Itô formula to this term, we rewrite $G$ as
\[G(\varphi, \hat{k}) = \left(I_3 + \sin(\varphi)Q(\hat{k}) + (1-\cos(\varphi))Q^2(\hat{k})\right)H_1(\hat{k})
\]
where $I_3$ is the $3 \times 3$ identity matrix, $Q$ is defined by (29), and where
\[H_1(\hat{k}) := \frac{1}{\sqrt{k_1^2 + k_2^2}} \begin{pmatrix} \frac{k_2}{k_1} & \frac{k_1^1}{k_1} \end{pmatrix} = \frac{1}{\sqrt{k_1^2 + k_2^2}} \begin{pmatrix} \frac{k_2^1}{k_1} & \frac{k_1^1}{k_1} \end{pmatrix},
\]
which is orthogonal to $\hat{k} = k/|k|$. In fact, $G(\varphi, \hat{k})$ corresponds to the rotation of $H_1(\hat{k}) \in \hat{k}^\perp$ with angle $\varphi$ and axis $\hat{k}$. This choice simplifies calculations. Now, note that
\[\int_{-1}^{1-\varepsilon} \rho(x, s)(s-1) ds \leq \int_{-1}^{1} \rho(x, s)(s-1) \leq \|a\| \frac{21-\eta_m/2}{1 - \eta_M/2},
\]
so that $E_{n,m}^{(7)}$ does not depend on $\varepsilon$. Applying the Itô formula, we obtain
\[
\sum_{n} \sum_{m} |E_{n,m}^{(7)}| \leq ht \left(1 + (1 + 2^7 \cdot 3^3)\sigma_{\infty,\varepsilon,0}^2 \|a\| \frac{21-\eta_m/2}{1 - \eta_M/2} \right) \\
\times \left( \|\nabla_x \lambda\|_\infty + \|\lambda\|_\infty + \|\nabla_x a\|_\infty \right) \left( \|D_k^1 \hat{b}\|_\infty + \|D_k^2 \hat{b}\|_\infty \right) =: hF_4(b)
\]
Setting finally $F_1 := F_2 + F_3 + F_4$ and gathering all previous results concludes the proof of Proposition 6.1. 

6.3.3 Step 3 and conclusion

We remark first that the error bound in Proposition 6.1 does not depend on the starting point $(x, \hat{k})$. Then, from this pointwise result, we find
\[
|\mathbb{E}_{\mu_0} [f(D_{h,\varepsilon}(T))] - \mathbb{E}_{\mu_0} [f(\hat{D}_c(T))]| \leq \int_{\mathbb{R}^3 \times \mathbb{S}^2} |\mathbb{E}_{x,\hat{k}} [f(D_{h,\varepsilon}(T))] - \mathbb{E}_{x,\hat{k}} [f(\hat{D}_c(T))]| \mu_0(dx, d\hat{k}),
\]

5.2.2 as an illustration. The goal is to simulate a real-valued random variable \( \alpha \) on the qualitative behavior of the solution.

In this section, we describe the stochastic collocation method, see e.g. [28], and consider the situation of Section 2.1 pp.16]. From either simulated or experimental data obtained e.g. from light propagation in biological tissues. This will require the development of appropriate inverse techniques.

We have derived an efficient MC method for the resolution of the RTE with non-integrable scattering kernels. It is based on a small jumps/large jumps decomposition that allows us to simulate the small jumps part at a low cost by solving a standard SDE. The large jumps are obtained by using the stochastic collocation technique with a candidate distribution function that captures the singular behavior of the kernel. We have moreover demonstrated the necessity to include the small jumps component in order to obtain a good accuracy at a manageable computational cost, and investigated practical situations in optical tomography and atmospheric turbulence where the singular RTE is of interest. We in particular highlighted the role of the singularity strength \( \alpha \) on the qualitative behavior of the solution.

Future investigations include the estimation of the scattering kernel, with an emphasis on the parameter \( \alpha \), from either simulated or experimental data obtained e.g. from light propagation in biological tissues. This problem is of practical interest in biomedical applications and will require the development of appropriate inverse techniques.

A Stochastic collocation

In this section, we describe the stochastic collocation method, see e.g. [28], and consider the situation of Section 5.2.2 as an illustration. The goal is to simulate a real-valued random variable \( W \) (for which direct simulation is not possible or too costly) from an auxiliary variable \( V \) that can be generated efficiently. In our context, we want to simulate \( W \) with probability density function (PDF)

\[
f_W(w) := \frac{a(\sqrt{2w})}{C_W w^{1+\alpha/2}} 1_{(\varepsilon,2)}(w),
\]

where \( C_W \) is a normalization constant. As already noticed in Section 3.4 a direct method is available when \( a \equiv 1 \). Therefore, we take \( V \) with PDF

\[
f_V(v) := \frac{1}{C_V v^{1+\alpha/2}} 1_{(\varepsilon,2)}(v),
\]

7 Conclusion

We have derived an efficient MC method for the resolution of the RTE with non-integrable scattering kernels. It is based on a small jumps/large jumps decomposition that allows us to simulate the small jumps part at a low cost by solving a standard SDE. The large jumps are obtained by using the stochastic collocation technique with a candidate distribution function that captures the singular behavior of the kernel. We have moreover demonstrated the necessity to include the small jumps component in order to obtain a good accuracy at a manageable computational cost, and investigated practical situations in optical tomography and atmospheric turbulence where the singular RTE is of interest. We in particular highlighted the role of the singularity strength \( \alpha \) on the qualitative behavior of the solution.

We conclude by applying the central limit theorem [16] together with the Portmanteau theorem [6, Theorem 2.1 pp.16].
that can be simulated with

\[ V = F_V^{-1}(U) = \varepsilon(1 - (1 - (\varepsilon/2)^{\alpha/2})U)^{-2/\alpha}, \]

where \( U \sim U(0, 1) \) and where \( F_V \) is the cumulative distribution function (CDF) of \( V \). The stochastic collocation method is based on the following three observations. First, we have \( F_V(V) \sim U(0, 1) \). Second, denoting by \( F_W \) the CDF of \( W \), we note that \( W \) can be (theoretically) simulated with

\[ F_W^{-1}(U) = F_V^{-1}(F_V(V)) =: G(V), \]

with \( G = F_W^{-1} \circ F_V \) and \( U = F_V(V) \). Last, we only need to approximate \( G \) and not \( F_W^{-1} \), and with a good candidate \( V \), \( G \) behaves better than \( F_W^{-1} \). In order to approximate \( G \), we use Gauss polynomial interpolation and only need to invert \( F_W \) at a small number of points.

In our example, \( V \) captures the "singular" behavior of \( W \), and is as a consequence a good candidate. The function \( G \) is then direct to approximate with just a few quadrature points for a reduced computational cost. Because \( G \) needs only to be approximated over \( (\varepsilon, 2) \), with known values at the extremes, we rather use a Gauss-Lobatto-Jacobi quadrature rule. In Figure 18, we illustrate the polynomial approximation of \( G \), with 5 and 10 interpolation points for \( \alpha = 5/3 \), \( \varepsilon = 0.01 \), and \( a(r) = \exp(-r^2/(2 \times 0.8^2)) \). Because of our choice for \( V \), one can observe that the overall behavior of the PDF \( f_W \) is well captured with just 5 quadrature points, even for strongly singular kernels with \( \alpha = 5/3 \). However, the fast decay of the function \( a \), which is the main source of error between \( F_W \) and \( F_V \), requires more quadrature points for an accurate approximation and 10 points seem sufficient.

![Figure 18: Illustration of the polynomial approximation to G with 5 (left picture) and 10 (right picture) interpolations points. We use the library Jacobi.jl to compute these quadrature points.](image)

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