Estimating stationary characteristic functions of stochastic systems via semidefinite programming

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Abstract—This paper proposes a methodology to estimate characteristic functions of stochastic differential equations that are defined over polynomials. For such systems, the time evolution of the characteristic function is governed by a partial differential equation; consequently, the stationary characteristic function can be obtained by solving an ordinary differential equation (ODE). However, except for a few special cases, the solution to the ODE consists of unknown coefficients. These coefficients are closely related with the stationary moments of the process, which could be estimated by utilizing the fact that the characteristic function is positive definite. The method is illustrated via examples.

I. Introduction

A plethora of systems in engineering, finance, and physical sciences exhibit stochastic dynamics [1]–[5]. These systems are mathematically characterized in terms of how their probability density functions (or characteristic functions) evolve with time. However, it is generally difficult to solve for these quantities other than a few special cases. One often resorts to stochastic simulations or other approximation schemes [6]–[10]. Here, we propose a method to estimate stationary characteristic functions for a class of stochastic systems. The state of this class evolves as per a stochastic differential equation driven by white noise and/or Lévy noise. For such systems, a partial differential equation (PDE) governing the time evolution of the characteristic function can be written [11], [12]. In stationary state, this PDE becomes an ordinary differential equation (ODE) whose order is determined by the degrees of polynomials in drift and diffusion terms. The solution to this ODE requires as many unknown coefficients as its order. Our method relies upon relating these coefficients to the stationary moments, and determining the bounds on them by exploiting the positive definiteness of the characteristic function.

Since the proposed method involves estimating some moments of a stochastic system, it is closely related with moment closure methods [13]-[15]. Specifically, recent works have shown that the positive semidefiniteness of the moments can be used to find exact lower and upper bounds on moments of a stochastic system [16]-[23]. Here, we show that such bounds can also be obtained by using the fact that the characteristic function is a positive definite function. Furthermore, the proposed method also utilizes the

moment bounds to approximate the stationary characteristic function. This work is also related with another category of works wherein the characteristic function is used to describe stochastic dynamical systems, including those driven by Poisson/Lévy noise [11], [12], [24]-[31]. However, these works are typically restricted to cases where the characteristic function can be exactly solved for. Our method approximates the characteristic function beyond these.

The rest of the paper is organized as follows. Section II provides background results on characteristic functions and their relevant properties. Section III describes generators of stochastic dynamical system and proposes our method to estimate moments and characteristic function in stationary state. Section IV illustrates the approach via examples, and finally section V concludes the paper.

Notation: The set of real numbers, and natural numbers are respectively denoted by \mathbb{R} and \mathbb{N} . The state of a stochastic process is denoted by $X \in \mathbb{R}$, with a specific value taken by it denoted by x. The squared root of -1 is denoted by the letter j. The expectation operator is denoted by \mathbb{E} . Characteristic function is denoted by $\varphi(\omega)$.

II. MATHEMATICAL PRELIMINARIES

In this section, we provide background results pertaining characteristic functions and their properties. The reader is referred to [32]-[35] for proofs and more details. For simplicity, we only consider one dimensional systems.

The characteristic function $\varphi: \mathbb{R} \to \mathbb{C}$ for a univariate random variable X is defined as

$$\varphi(\omega) := \mathbb{E}\left(e^{j\omega X}\right) = \int_{-\infty}^{\infty} e^{j\omega x} dF(x),\tag{1}$$

where F is the distribution function of X. For any random variable, the characteristic function always exists and it uniquely determines the distribution. If F(x) has a density h(x) then $\varphi(\omega)$ can be written in terms of h

$$\varphi(\omega) = \int_{-\infty}^{\infty} e^{j\omega x} h(x) dx. \tag{2}$$

Given $\varphi(\omega)$, the distribution function F(x) and/or the density h(x) can be obtained via inversion.

As an immediate consequence of its definition, the characteristic function of a random variable has the following properties

- 1) $\varphi(0) = 1$.
- 2) $|\varphi(\omega)| \le 1$. 3) $\varphi(-\omega) = \overline{\varphi(\omega)}$.

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4) If the random variable has a finite m^{th} order moment, then it is given by

$$\mathbb{E}(X^m) = (-j)^m \frac{d^m \varphi}{d\omega^m}|_{\omega=0}.$$
 (3)

Another important property of a characteristic function that we particularly use in this work is that it is a positive definite function. The following theorem of Bochner forms basis of our analysis.

Theorem 1 (Bochner): A continuous complex-valued function φ on \mathbb{R} is a characteristic function if and only if $\varphi(0) = 1$ and φ is positive definite [33], [34].

The positive definiteness of a function required by this theorem is defined as follows. A complex valued function $\varphi:\mathbb{R}\to\mathbb{C}$ is said to be positive definite if the inequality

$$\sum_{i,j=1}^{q} \varphi(\omega_i - \omega_j) c_i \overline{c_j} \ge 0 \tag{4}$$

holds for every positive integer q, for all $\omega_1,\ldots,\omega_q\in\mathbb{R}$, and for all $c_1,\ldots,c_q\in\mathbb{C}$. In other words, φ is positive definite if and only if the matrix

$$M = \begin{bmatrix} \varphi(0) & \varphi(\omega_1 - \omega_2) & \dots & \varphi(\omega_1 - \omega_q) \\ \varphi(\omega_2 - \omega_1) & \varphi(0) & \dots & \varphi(\omega_2 - \omega_q) \\ \vdots & \vdots & \ddots & \vdots \\ \varphi(\omega_q - \omega_1) & \varphi(\omega_q - \omega_2) & \dots & \varphi(0) \end{bmatrix} \succeq 0.$$
(5)

for an arbitrary choice of $q \in \mathbb{N}$ and $\omega_1, \dots, \omega_q \in \mathbb{R}$. Here $\succeq 0$ denotes the positive semidefiniteness.

As a consequence of this definition, a variety of properties of the characteristic function (including properties 2 and 3 above) can be established by choosing some test points and enforcing $M\succeq 0$. For example, consider the case q=1. Then, we must have that

$$\varphi(0) \ge 0. \tag{6}$$

Likewise, for q = 2, we should have

$$\begin{bmatrix} \varphi(0) & \varphi(\omega_1 - \omega_2) \\ \varphi(\omega_2 - \omega_1) & \varphi(0) \end{bmatrix} \succeq 0. \tag{7}$$

Without loss of generality, we can assume $\omega_1 = 0$ and $\omega_2 = \omega \in \mathbb{R}$. Then, we should have $\varphi(-\omega) = \overline{\varphi(\omega)}$ and $|\varphi(\omega)|^2 \le \varphi(0)^2$.

III. STATIONARY CHARACTERISTIC FUNCTION OF A STOCHASTIC PROCESS

In this section, we describe the generator of a stochastic process and get a PDE for time evolution of the characteristic function. Then, we use this PDE to obtain an ODE for the stationary characteristic function and discuss its solution.

A. Stochastic dynamical system and its generator

We consider stochastic differential equations driven by Lévy noise, known as *Itô-Lévy* processes. In this paper, we will restrict ourselves to Itô-Lévy processes of the form:

$$dX(t) = f(X(t))dt + g(X(t))dW(t) + \int_{R} y\tilde{N}(dt, dy).$$
 (8)

Here $X \in \mathbb{R}$ denotes the state, W is the Wiener process, \tilde{N} is called a *compensated Poisson measure*, and $f: \mathbb{R} \to \mathbb{R}$, $g: \mathbb{R} \to \mathbb{R}$ characterize the dynamics. In order for (8) to be well-defined we must assume that the driving Lévy process has finite variance. See [36], [37] for more details on the formalism.

The jump process generalizes the jumps of a Poisson process. In a Poisson process, all of the jumps have value 1, and N(t) simply counts the number of jumps. In a Lévy process, the jumps can take on arbitrary values. The *Poisson random measure*, $N(t,\cdot)$, is a measure-valued stochastic process such that for a Borel set S, N(t,S) is a Poisson process that counts the number of jumps that took values in the set S. The intensity of the Poisson random measure given by the *Lévy measure*, ν :

$$\mathbb{E}[N(t,S)] = t\nu(S)$$

The compensated Poisson measure from (8) is the measure-valued stochastic process defined by $\tilde{N}(t,S) = N(t,S) - t\nu(S)$.

Let ψ be a twice continuously differentiable function. A standard result in stochastic differential equations [36] shows that the dynamics of $\mathbb{E}[\psi(X(t))]$ are given by:

$$\frac{d}{dt}\mathbb{E}[\psi(X(t))] = \mathbb{E}[L\psi(X(t))] \tag{9}$$

where $L\psi$ is called the *generator* of the process. The generator is defined by

$$L\psi(X) = \frac{\partial \psi(X)}{\partial X} f(X) + \frac{1}{2} \frac{\partial^2 \psi(X)}{\partial X^2} g(X)^2 + \int_{\mathbb{R}} \left(\psi(X+y) - \psi(X) - \frac{\partial \psi(X)}{\partial X} y \right) \nu(dy). \quad (10)$$

In the following, we use the generator to find a PDE that governs the evolution of the characteristic function. We will see that for the commonly-studied Lévy processes, formulas exist to enable explicit calculation of the required integral.

B. Characteristic function of the process

Consider the stochastic dynamics defined in (8). We restrict ourselves to the cases for which the functions f, and g are polynomials in the state X. Our goal is to compute the time evolution of

$$\psi(X(t)) = e^{j\omega X(t)}. (11)$$

The following theorem shows how a partial differential equation governing the evolution of the characteristic function can be obtained for (8). We wish to point out that it is presented here for a formal statement, and it has been used in some form or other in several works, e.g., see [12].

Theorem 2: Consider a one dimensional stochastic process defined in (8). Let f and g^2 be finite polynomials of degrees $d_f \in \mathbb{N}$ and $d_g \in \mathbb{N}$ respectively. Assuming that a stationary distribution with finite moments of order $n = \max\{d_f, d_g\}$ exists, the characteristic function of

the stationary distribution satisfies the following ordinary differential equation of order n:

$$j\omega \sum_{l=0}^{d_f} a_{f_l} j^{-l} \frac{\partial^l \varphi}{\partial \omega^l} + \frac{1}{2} (j\omega)^2 \sum_{l=0}^{d_g} a_{g_l} j^{-l} \frac{\partial^l \varphi}{\partial \omega^l} + \eta(\omega) \varphi = 0,$$
(12)

where $\eta(\omega) = \int_{\mathbb{R}} \left(e^{j\omega y} - 1 - j\omega y \right) \nu(dy)$.

Proof: Since we assume that f, and g^2 are polynomials, without loss of generality we can take their forms to be

$$f(X) = \sum_{l=0}^{d_f} a_{f_l} X^l, \quad g^2(X) = \sum_{l=0}^{d_g} a_{g_l} X^l, \quad (13)$$

where $a_{f_l} \in \mathbb{R}$, and $a_{g_l} \in \mathbb{R}$ are coefficients. Taking $\psi(X) = e^{j\omega X}, \ \omega \in \mathbb{R}$, we have that

$$\frac{\partial^m \psi}{\partial X^m} = (j\omega)^m \psi, \quad \frac{\partial^m \psi}{\partial \omega^m} = (jX)^m \psi \quad \forall m \in \mathbb{N} \quad (14)$$

Using these, we can write $L\psi(X)$ as

$$L\psi(X) = j\omega\psi(X) \sum_{l=0}^{d_f} a_{f_l} X^l + \frac{1}{2} (j\omega)^2 \psi(X) \sum_{l=0}^{d_g} a_{g_l} X^l + \psi(X) \int_{\mathbb{R}} \left(e^{j\omega y} - 1 - j\omega y \right) \nu(dy).$$
 (15)

Using (14), the terms $X^l \psi$ can be replaced by $(j)^{-l} \frac{\partial^l \psi}{\partial \omega^l}$. This yields

$$L\psi(X) = j\omega \sum_{l=0}^{d_f} a_{f_l} j^{-l} \frac{\partial^l \psi}{\partial \omega^l} + \frac{1}{2} (j\omega)^2 \sum_{l=0}^{d_g} a_{g_l} j^{-l} \frac{\partial^l \psi}{\partial \omega^l} + \psi(X) \int_{\mathbb{R}} \left(e^{j\omega y} - 1 - j\omega y \right) \nu(dy). \quad (16)$$

Taking expectation, we get a partial differential equation in characteristic function $\varphi=\mathbb{E}(\psi(X))$

$$\frac{\partial \varphi}{\partial t} = j\omega \sum_{l=0}^{d_f} a_{f_l} j^{-l} \frac{\partial^l \varphi}{\partial \omega^l} + \frac{1}{2} (j\omega)^2 \sum_{l=0}^{d_g} a_{g_l} j^{-l} \frac{\partial^l \varphi}{\partial \omega^l} + \eta(\omega)\varphi, \quad (17)$$

where we used $\eta(\omega)$ to denote the integral $\int_{\mathbb{R}} \left(e^{j\omega y}-1-j\omega y\right) \nu(dy)$. If the stationary distribution with finite moments of order n exists, then we must have that $\frac{\partial \varphi}{\partial t}=0$. This results in the ordinary differential equation (12). The degree of this ODE is $n=\max\{d_f,d_g\}$. It is worth pointing out that there are pathological cases wherein expectation of a test function may not be appropriately defined (see [38] for some examples), the finite moments requirement allows us to avoid them [17].

Remark 1: The function $\eta(\omega)$ is known in as the *characteristic exponent*, [37]. A special case of the Lévy-Khintchine formula shows that the compensated Poisson process has characteristic function given by:

$$\mathbb{E}\left[e^{j\omega\int_0^t y\tilde{N}(dt,dy)}\right] = e^{t\eta(\omega)}.$$
 (18)

Thus, the statistics of the jump process $\int_0^t y \tilde{N}(dt, dy)$ are entirely determined by $\eta(\omega)$.

For commonly-studied Lévy processes, the characteristic exponent has an explicit formula. For example, the gamma process $\tau(t)$ has Lévy measure

$$\nu(dy) = ay^{-1}e^{-by}dy,\tag{19}$$

where a and b are positive parameters. It can be shown that

$$\mathbb{E}[e^{j\omega\tau(t)}] = e^{-ta\log\left(1 - j\frac{\omega}{b}\right)} \quad \text{ and } \quad \mathbb{E}[\tau(t)] = t\frac{a}{b} \quad (20)$$

It follows that the compensated Gamma process, $\tau(t) - \mathbb{E}[\tau(t)]$, has characteristic exponent given by $-a\log\left(1-j\frac{\omega}{b}\right) -j\omega a/b$.

A more complex example, which will be studied in an example below is the variance-gamma process: $\sigma W(\tau(t))$. This process is formed by composing a standard Brownian motion W(t) with a gamma process. In this case, it can be shown that

$$\mathbb{E}[e^{j\omega W(\tau(t))}] = e^{-ta\log\left(1 + \frac{(\sigma\omega)^2}{2b}\right)} \quad \text{and} \quad \mathbb{E}[W(\tau(t))] = 0. \tag{21}$$

It follows that compensated variance-gamma process is simply the variance-gamma process. Furthermore, the characteristic exponent is given by

$$\eta(\omega) = -a \log \left(1 + \frac{(\sigma \omega)^2}{2b} \right),$$
(22)

Next, we discuss the solution of the ODE for the stationary characteristic function.

C. Solving for stationary characteristic function

The ODE in (12) cannot be solved analytically except for a handful of cases. However, it can be solved via numerical techniques. Either way it would require n initial/intermediate/boundary values to find the solution. Other than the usual $\varphi(0)=1$, previous works have either utilized prior knowledge about the system (e.g., the distribution is symmetric), or used $\lim_{|\omega|\to\infty}\varphi(\omega)=0$ and $\lim_{|\omega|\to\infty}\frac{\partial^l\varphi(\omega)}{\partial\omega^l}=0$ for some l [12]. In practice these are hard to incorporate in a solution. Furthermore, if one is interested only in stationary moments, then solution of $\varphi(\omega)$ only in neighborhood of zero is sufficient.

We propose a different approach to compute both the moments and the characteristic function. This approach relies on two ideas. First being the fact that the characteristic function is related with the moments as

$$\frac{\partial^l \varphi}{\partial \omega^l} |_{\omega=0} = j^l \mu_l, \ l = \{1, \dots, n-1\}, \tag{23}$$

where $\mu_l \in \mathbb{R}$ represents the l^{th} order moment. Thus, the moments are natural quantities to be used in computing the characteristic function. Second idea is to utilize the Bochner's theorem to estimate the moments μ_l . In particular, we can use $\varphi(0)=1$ and positive semidefinite property of the matrix M in (5). Using these, a semidefinite program can be formulated that gives lower and upper bounds on μ_l as stated in the theorem below.

Theorem 3: Consider the one dimensional system defined by (8). Assuming that f and g^2 are polynomials of degree d_f and d_g , a lower bound on a moment μ_l can be obtained via the semidefinite program

$$\min \mu_k \tag{24a}$$

$$j\omega \sum_{l=0}^{d_f} a_{f_l} j^{-l} \frac{\partial^l \varphi}{\partial \omega^l} + \frac{1}{2} (j\omega)^2 \sum_{l=0}^{d_g} a_{g_l} j^{-l} \frac{\partial^l \varphi}{\partial \omega^l} + \eta(\omega) \varphi = 0,$$
(24b)

$$\frac{\partial^{l} \varphi}{\partial \omega^{l}} |_{\omega=0} = j^{l} \mu_{l}, \quad l = 1, 2, \dots n, \tag{24c}$$

$$\varphi(0) = 1, (24d)$$

$$M \succeq 0.$$
 (24e)

Here $k = \{1, ..., n\}$ with $n = \max\{d_f, d_g\}$ and M is defined as in (5). Further, the minimum value obtained by the program increases as size of M is increased by including more test points.

Proof: Since f and g^2 are assumed to be polynomials, Theorem 2 implies that the characteristic function $\varphi(\omega)$ satisfies the ODE of order n given by (24b). The moments are related with the derivatives of the characteristic function by the linear constraints in (24c). The constraint in (24d) and (24e) are a consequence of the Bochner's theorem. Since the objective function is linear in decision variables μ_k and the constraints are either equality or semidefinite constraints, the optimization problem is a semidefinite program [39].

Now suppose that the size of M is increased by including more test points $\omega_1, \ldots, \omega_q$. This corresponds to adding more constraints in the program, and the solution cannot get worse by doing so.

The upper bound on μ_k can be found by minimizing $-\mu_k$. Note that we can choose any test points $\omega_1, \omega_2, \ldots, \omega_q \in \mathbb{R}$ in order to generate the matrix M. For sake of simplicity, we will choose uniformly spaced values on the real-line.

The above semidefinite program can be used to compute lower and upper bounds on each of the moments μ_k . These values can be then used to determine the solution to the ODE for characteristic function and thereby finding an approximation of the characteristic function. If the lower and upper bounds on each of the moments are reasonably close, then the approximate characteristic function would be quite close to the true characteristic function. It can be further used to compute the stationary probability density via inversion.

Remark 2: The formulation in (24) may not always be convenient, specially if (24b) does not have an analytical solution. In this case, we can interpret (24) as an optimal control problem for a linear time varying system. Specifically, consider a state vector $z = \left[\varphi \quad \frac{\partial \varphi}{\partial \omega} \quad \cdots \quad \frac{\partial^{n-1} \varphi}{\partial \omega^{n-1}} \right]^{\top}$. Then, the differential equation describing the characteristic function becomes a linear time (in ω -space) varying system

$$\frac{dz}{d\omega} = \mathcal{P}(\omega)z(\omega),\tag{25}$$

for an appropriately defined matrix \mathcal{P} . In this setup, the objective would be to optimize the elements of z(0) subject

to the linear matrix inequality (24e). The linear system in (25) can be appropriately discretized, and the decision variables of the optimization problem are values of the the state vector z at those discrete points. The solution to this semidefinite program incurs numerical errors due to discretization.

Remark 3: The semidefinite program in (24) can be used to find bounds on first n moments where n is the degree of the ODE (24b). If one is interested in computing the higher order moments, the approximate characteristic function can be differentiated and computed at $\omega=0$. By doing so, bounds on the higher order moments can also be computed. Alternatively, for systems with finite moments, one can compute the bounds on first n moments via the proposed method, and then use the fact that stationary moments are related via a linear system of equations given by [18]

$$A\mathcal{X} + B\overline{\mathcal{X}} = 0. {26}$$

Here \mathcal{X} is collection of moments up to some order and $\overline{\mathcal{X}}$ contains moments of order higher than those in \mathcal{X} . The number of elements of $\overline{\mathcal{X}}$ is as many as the degree of nonlinearity in the system given by n. While the usual moment closure methods estimate elements of $\overline{\mathcal{X}}$ in terms of those of \mathcal{X} , we simply supplement (26) with lower and upper bounds on n moments and thereby compute bounds on all other moments in (26).

In the next section, we illustrate the proposed method on some simple examples and verify its performance.

IV. EXAMPLES

Example 1 (Stochastic Logistic Model): Consider the following modified stochastic logistic growth model

$$dX = (1 + X - 0.1X^{2})dt + \sqrt{2}XdW. \tag{27}$$

Without the constant term in the drift, this model is widely-used in modeling population growth [1]. We added the constant term so that the trivial solution X=0 is ruled out.

Using (17), the characteristic function evolves as per

$$\frac{\partial \varphi}{\partial t} = j\omega\varphi + \omega \frac{\partial \varphi}{\partial \omega} + j0.1\omega \frac{\partial^2 \varphi}{\partial \omega^2} + \omega^2 \frac{\partial^2 \varphi}{\partial \omega^2}.$$
 (28)

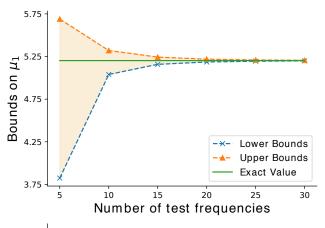
Therefore, the stationary characteristic function is the solution to the following differential equation

$$j\omega\varphi + \omega\frac{\partial\varphi}{\partial\omega} + (j0.1\omega + \omega^2)\frac{\partial^2\varphi}{\partial\omega^2} = 0.$$
 (29)

It can be shown that the above ODE has the following generalized solution

$$\varphi(\omega) = \frac{c_1}{\sqrt{5}} I_0 \left(2\sqrt{0.1 - j\omega} \right) + \frac{c_2}{\sqrt{5}} K_0 \left(2\sqrt{0.1 - j\omega} \right), \tag{30}$$

where I and K denote the modified Bessel functions of first and second kinds, and c_1 , c_2 are unknown coefficients. As expected, the number of unknown coefficients is same as the order of nonlinearity in the dynamics.



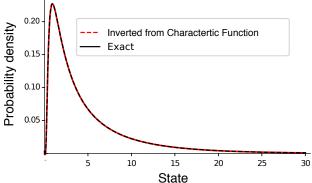


Fig. 1. Comparison of estimated mean and characteristic function with exact values. Top: The mean μ_1 is estimated via semidefinite program in (24). The matrix M is constructed by taking frequencies spaced by one (e.g., $\omega_1=1,\omega_2=2,\ldots$). Increasing the number of test frequencies leads to a better estimation of the moment as both lower and upper bounds converge. Estimation of mean via SDP. Bottom: The characteristic function is estimated by using the middle value of lower and upper bounds for μ_1 obtained from 30 test frequencies in (30)–(31). The estimated probability densify function is obtained by numerically inverting the characteristic function. Notably, this particular system is exactly solvable, and the exact probability densify function is provided for comparison purpose, showing an excellent match with its estimated counterpart.

To determine the coefficients, we can use the fact that $\varphi(0)=1$ and $\frac{\partial \varphi}{\partial \omega}|_{\omega=0}=j\mu_1$, where $\mu_1\in\mathbb{R}$ is the mean that is to be determined. This results in

$$\frac{c_1}{\sqrt{5}}I_0\left(2\sqrt{0.1}\right) + \frac{c_2}{\sqrt{5}}K_0\left(2\sqrt{0.1}\right) = 1 \tag{31a}$$

$$-j\sqrt{2}c_1I_1\left(2\sqrt{0.1}\right) + j\sqrt{2}c_2K_1\left(2\sqrt{0.1}\right) = j\mu_1$$
 (31b)

Using these, the stationary characteristic function can be written in terms of only one unknown μ_1 . Now, we can compute bounds on μ_1 using the semidefinite program as in (24). By choosing uniformly spaced values of $\omega_1, \omega_2, \ldots, \omega_n$, we computed the maximum and minimum allowable values μ_1 . We also find that increasing the size of the program by choosing more test points improves both lower and upper bounds. Taking 30 test points at $\omega_1 = 1, \ldots \omega_{30} = 30$, we get $5.2024 \le \mu_1 \le 5.2025$ (see Fig. 1, Top). This result is in excellent agreement with Monte Carlo simulations which yield an value of 5.2 for 10000 simulations.

Using the value of μ_1 obtained here, we can use the characteristic function to reconstruct the probability density

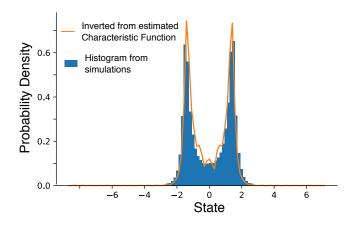


Fig. 2. Comparison of probability density functions obtained from the proposed method and simulations. The histogram is obtained via 1000 simulations of the process, whereas the estimated probability density is inverted from the estimated characteristic function by solving the semidefinite program in (24). The optimization was performed by casting the problem as an optimal control problem described in Remark 2, with evenly spaced 80 test frequencies.

function of the stationary distribution (see Fig. 1, Bottom). As mentioned earlier, the bounds on μ_1 can be used to estimate bounds on higher order moments as well (results not shown here).

Example 2 (Variance Gamma Process): Consider the following process

$$dX = (-5X^3 + 10X)dt + dW(\tau(t)), \tag{32}$$

where $W(\tau(t))$ is the variance-gamma process described above. Recall, in particular, $W(\tau(t))$ has zero mean, and so the compensated variance-gamma process is the same as the original variance-gamma process. Furthermore, recall the expression for the characteristic exponent, $\eta(\omega)$, from (22).

Using (17), the characteristic function of this process is governed by

$$\frac{\partial \varphi}{\partial t} = 5\omega^3 \frac{\partial^3 \varphi}{\partial \omega^3} + 10\omega \frac{\partial \varphi}{\partial \omega} + \eta(\omega)\varphi \tag{33}$$

As before, this ODE solution will have three unknown coefficients, out of which two could be readily computed by using the facts that $\varphi(0)=1$ and the stationary mean value of X is zero due to symmetry of the process. Numerically solving the optimization problem that minimizes μ_2 (the second moment), and inverting the solution yields quite accurate estimate of the stationary probability density function (see Fig. 2).

The optimization was performed by discretizing the optimal control problem described in Remark 2 via the trapezoidal rule, and then directly solving the corresponding SDP. Here, the decision variables were the state variables $z(\omega_i)$, where ω_i are the discretization frequencies. One subtlety is that in order to express the matrix M in terms of decision variables, the frequencies, ω_i , must be evenly spaced.

V. CONCLUSION

In this paper, we proposed a method to estimate stationary characteristic function for a class of stochastic dynamical system driven by both white and Lévy noise. The characteristic function for these stochastic systems is governed by an ODE. The method relies upon restricting the solutions of the ODE to positive definite functions and casts the problem as a semidefinite program. In future work, we would extend this framework to multidimensional cases. It would also be interesting to explore whether transient solutions of the characteristic function (described via PDE) can be obtained via a similar method.

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REFERENCES

- E. Allen, Modeling with Itô stochastic differential equations, vol. 22. Springer Science & Business Media, 2007.
- [2] R. Lande, S. Engen, and B.-E. Saether, Stochastic population dynamics in ecology and conservation. Oxford University Press on Demand, 2003
- [3] A. G. Malliaris, Stochastic methods in economics and finance, vol. 17. North-Holland, 1982.
- [4] C. Gardiner, "Handbook of stochastic methods for physics, chemistry and the natural sciences," *Applied Optics*, vol. 25, p. 3145, 1986.
- [5] B. Øksendal, Stochastic differential equations. Springer, 2003.
- [6] H. Risken and T. Frank, The Fokker-Planck Equation: Methods of Solution and Applications. Springer Series in Synergetics, 1996.
- [7] L. Socha, Linearization Methods for Stochastic Dynamic Systems. Lecture Notes in Physics 730, Springer-Verlag, Berlin Heidelberg, 2008
- [8] L. Socha, "Linearization in analysis of nonlinear stochastic systems: Recent results–part i: Theory," *Applied Mechanics Reviews*, vol. 58, no. 3, pp. 178–205, 2005.
- [9] J. P. Hespanha and A. Singh, "Stochastic models for chemically reacting systems using polynomial stochastic hybrid systems," *International Journal of robust and nonlinear control*, vol. 15, pp. 669–689, 2005.
- [10] J. P. Hespanha, "Modelling and analysis of stochastic hybrid systems," IEE Proceedings Control Theory And Applications, vol. 153, no. 5, p. 520, 2006.
- [11] M. Grigoriu, "Characteristic function equations for the state of dynamic systems with gaussian, poisson, and levy white noise," *Probabilistic Engineering Mechanics*, vol. 19, pp. 449–461, 2004.
- [12] M. Grigoriu, Stochastic calculus: applications in science and engineering. Springer Science & Business Media, 2013.
- [13] C. Kuehn, Moment Closure—A Brief Review. Understanding Complex Systems, Springer, 2016.
- [14] M. Soltani, C. A. Vargas-Garcia, and A. Singh, "Conditional moment closure schemes for studying stochastic dynamics of genetic circuits," *Biomedical Circuits and Systems, IEEE Transactions on*, vol. 9, no. 4, pp. 518–526, 2015.
- [15] A. Singh and J. P. Hespanha, "Approximate moment dynamics for chemically reacting systems," *IEEE Transactions on Automatic Con*trol, vol. 56, no. 2, pp. 414–418, 2011.
- [16] A. Lamperski, K. R. Ghusinga, and A. Singh, "Stochastic optimal control using semidefinite programming for moment dynamics," in *De*cision and Control (CDC), 2016 IEEE 55th Conference on, pp. 1990– 1995, IEEE, 2016.
- [19] K. R. Ghusinga, C. A. Vargas-Garcia, A. Lamperski, and A. Singh, "Exact lower and upper bounds on stationary moments in stochastic biochemical systems," *Physical Biology*, vol. 14, 2017.

- [17] J. Kuntz, M. Ottobre, G.-B. Stan, and M. Barahona, "Bounding stationary averages of polynomial diffusions via semidefinite programming," *SIAM Journal on Scientific Computing*, vol. 38, no. 6, pp. A3891–A3920, 2016.
- [18] A. Lamperski, K. R. Ghusinga, and A. Singh, "Analysis and control of stochastic systems using semidefinite programming over moments," arXiv preprint arXiv:1702.00422, 2017.
- [20] J. Kuntz, P. Thomas, G.-B. Stan, and M. Barahona, "Rigorous bounds on the stationary distributions of the chemical master equation via mathematical programming," arXiv preprint arXiv:1702.05468, 2017.
- [21] Y. Sakurai and Y. Hori, "Optimization-based synthesis of stochastic biocircuits with statistical specifications," *Journal of The Royal Society Interface*, vol. 15, no. 138, p. 20170709, 2018.
- [22] G. R. Dowdy and P. I. Barton, "Bounds on stochastic chemical kinetic systems at steady state," *The Journal of Chemical Physics*, vol. 148, no. 8, p. 084106, 2018.
- [23] K. R. Ghusinga, A. Lamperski, and A. Singh, "Moment analysis of stochastic hybrid systems using semidefinite programming," arXiv preprint arXiv:1802.00376, 2018.
- [24] T. Epps, "Characteristic functions and their empirical counterparts: geometrical interpretations and applications to statistical inference," *The American Statistician*, vol. 47, pp. 33–38, 1993.
- [25] M. Grigoriu, "A partial differential equation for the characteristic function of the response of non-linear systems to additive poisson white noise," *Journal of Sound and vibration*, vol. 198, pp. 193–202, 1996.
- [26] S. F. Wojtkiewicz, E. A. Johnson, L. A. Bergman, M. Grigoriu, and B. F. Spencer, "Response of stochastic dynamical systems driven by additive gaussian and poisson white noise: Solution of a forward generalized kolmogorov equation by a spectral finite difference method," *Computer methods in applied mechanics and engineering*, vol. 168, pp. 73–89, 1999.
- [27] A. Chechkin, V. Gonchar, J. Klafter, R. Metzler, and L. Tanatarov, "Stationary states of non-linear oscillators driven by lévy noise," *Chemical Physics*, vol. 284, pp. 233–251, 2002.
- [28] B. Dybiec, E. Gudowska-Nowak, and I. Sokolov, "Stationary states in langevin dynamics under asymmetric lévy noises," *Physical Review E*, vol. 76, p. 041122, 2007.
- [29] G. Samorodnitsky and M. Grigoriu, "Characteristic function for the stationary state of a one-dimensional dynamical system with lévy noise," *Theoretical and Mathematical Physics*, vol. 150, pp. 332–346, 2007.
- [30] S. Denisov, W. Horsthemke, and P. Hänggi, "Generalized fokkerplanck equation: Derivation and exact solutions," *The European Physical Journal B-Condensed Matter and Complex Systems*, vol. 68, pp. 567–575, 2009.
- [31] G. Cottone, "Statistics of nonlinear stochastic dynamical systems under lévy noises by a convolution quadrature approach," *Journal of Physics A: Mathematical and Theoretical*, vol. 44, p. 185001, 2011.
- [32] W. Feller, An introduction to probability theory and its applications: volume II. John Wiley & Sons New York, 1968.
- [33] N. G. Ushakov, Selected topics in characteristic functions. Walter de Gruyter, 1999.
- [34] T. M. Bisgaard and Z. Sasvári, Characteristic Functions and Moment Sequences: Positive Definiteness in Probability. Nova Publishers, 2000.
- [35] Z. Sasvári, Multivariate characteristic and correlation functions. Walter de Gruyter, 2013.
- [36] B. Øksendal and A. Sulem, Applied Stochastic Control of Jump Diffusions. Springer, 3 ed., 2009.
- [37] D. Applebaum, Lévy processes and stochastic calculus. Cambridge university press, 2009.
- [38] P. W. Glynn and A. Zeevi, "Bounding stationary expectations of markov processes," in *Markov processes and related topics: a Festschrift for Thomas G. Kurtz*, pp. 195–214, Institute of Mathematical Statistics, 2008.
- [39] S. Boyd and L. Vandenberghe, *Convex optimization*. Cambridge university press, 2004.