Principal axes for stochastic dynamics

V.V. Vasconcelos,1 F. Raichsel,2 M. Haase,3 J. Peinke,4 M. Wächter,4 P.G. Lind,1,2 and D. Kleinhan5,6

1Physics Department, Faculty of Sciences, University of Lisbon, 1649-003 Lisbon, Portugal
2Center for Theoretical and Computational Physics, University of Lisbon, Av. Prof. Gama Pinto 2, 1649-003 Lisbon, Portugal
3Institute for High Performance Computing, University of Stuttgart, Nobelstr. 19, DE-70569 Stuttgart, Germany
4Institute of Physics, University of Oldenburg, DE-26111 Oldenburg, Germany
5Institute for Marine Ecology, University of Gothenburg, Box 461, SE-405 30 Göteborg, Sweden
6Institute of Theoretical Physics, University of Münster, DE-48149 Münster, Germany

We introduce a general procedure for directly ascertaining how many independent stochastic sources exist in a complex system modeled through a set of coupled Langevin equations of arbitrary dimension. The procedure is based on the computation of the eigenvalues and the corresponding eigenvectors of local diffusion matrices. We demonstrate our algorithm by applying it to two examples of systems showing Hopf-bifurcation. We argue that computing the eigenvectors associated to the eigenvalues of the diffusion matrix at local mesh points in the phase space enables one to define vector fields of stochastic eigendirections. In particular, the eigenvector associated to the lowest eigenvalue defines the path of minimum stochastic forcing in phase space, and a transform to a new coordinate system aligned with the eigenvectors can increase the predictability of the system.

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I. INTRODUCTION

When dealing with measurements on complex systems it is typically difficult to find the optimal set of variables to describe their evolution, which provides the best opportunities for understanding and predicting the system’s behavior. Recently a framework for analyzing measurements on stochastic systems was introduced [1, 2]. This framework is based on the assumption that the measured properties evolve according to a deterministic drift and stochastic fluctuations due to the interactions with the internal degrees of freedom or the environment. These fluctuations are globally ruled by some specific stochastic forcing inherent to the signal evolution itself, and therefore cannot be separated from the measured variables. Whether they are due to quantum uncertainties or a limited knowledge of the state of the system, they are fundamental for a complete description of a complex system. Several works in this scope have already shown the advantage of this approach, ranging from the description of turbulent flows [1] and climate indices [3] to the evolution of stock markets [4] and oil prices [5], just to mention some. In the course of time several improvements to the method were proposed concerning its robustness with respect to finite sampling effects and measurement noise [6,11].

However, a natural question arises when dealing simultaneously with several variables: is it possible to find non-trivial functions of measured properties for which the stochastic fluctuation can be neglected? This would mean that it should be possible to decrease the number of stochastic variables needed to describe the system.

In this paper we will follow this question up in detail and show, that by accessing the eigenvalues of the diffusion matrices comprehending the measured properties it is possible to derive a path in phase space through which the deterministic contribution is enhanced. As a direct application, our procedure allows for the determination of the number of independent sources of stochastic forcing in a system described by an arbitrarily large number of properties.

A pictorial example is as follows. Consider a professional archer trying to aim at a target in a succession of shots. If the archer holds the bow without any support, one expects that the set of trials is distributed around the center according to a radially symmetric Gaussian. Hence deviations have the same amplitudes in all directions. However, if the archer lies on the floor, the vertical direction will be more confined than the horizontal. In this case we expect a smaller variance of deviations in the vertical direction than in the horizontal one. If the archer lies on an inclined plane the most confined and less confined directions will have a certain slope related to the plane inclination. Regarding the variance in space as a particular representation of the diffusion matrix, we can now think in more complex systems where at each point in phase space fluctuations can be defined through their local diffusion matrices. The study of its eigenvalues and eigenvectors is the scope of the present paper.

We start in Sec. [II] by defining our system of variables mathematically and by describing the standard approach to estimate its dynamics from measured data as introduced in Ref. [1, 2], where a particular emphasis will be given to the diffusive terms. Then the general eigenvalue problem is introduced in the role of the local eigenvalues of diffusion matrices are discussed, which exhibits the framework for the present work. We show that at each point of phase space the eigenvectors of the diffusion matrix are tangent to new coordinate lines, one of them corresponding to the lowest eigenvalue, thus indicating the direction towards which fluctuations are - at least partially - suppressed. In Sec. [III] we demonstrate our approach on two examples from Hopf-bifurcation systems with stochastic forcing. Section [IV] closes the paper with discussion and conclusions.
II. DEFINING STOCHASTIC EIGENDIRECTIONS

We consider an $N$-dimensional Langevin process $X = (X_1(t), \ldots, X_N(t))$ whose probability density functions (PDFs) $f(X,t)$ evolve according to the Fokker-Planck equation (FPE) [12,13]

$$\frac{\partial f(X,t)}{\partial t} = -\sum_{i=1}^{N} \frac{\partial}{\partial x_i} \left[D_i^{(1)}(X)f(X,t)\right] + \sum_{i=1}^{N} \sum_{j=1}^{N} \frac{\partial^2}{\partial x_i \partial x_j} \left[D_{ij}^{(2)}(X)f(X,t)\right] \tag{1}$$

The functions $D_i^{(1)}$ and $D_{ij}^{(2)}$ are called the Kramers-Moyal or drift and diffusion coefficients and are defined as

$$D^{(k)}(X) = \lim_{\Delta t \to 0} \frac{1}{\Delta t} \frac{M^{(k)}(X, \Delta t)}{k!}, \quad (2)$$

where $M^{(k)}$ are the first and second conditional moments ($k = 1, 2$). Here we assume that the underlying process is stationary and therefore both drift and diffusion coefficients do not explicitly depend on time $t$. Conditional moments can be directly derived from the measured data as [2,7]:

$$M_i^{(1)}(X, \Delta t) = \langle Y_i(t+\Delta t) - Y_i(t) \rangle | Y(t) = X$$

$$M_{ij}^{(2)}(X, \Delta t) = \langle (Y_i(t+\Delta t) - Y_i(t))(Y_j(t+\Delta t) - Y_j(t)) \rangle | Y(t) = X \tag{3}$$

Here $Y(t) = (Y_1(t), \ldots, Y_N(t))$ exhibits the $N$-dimensional vector of measured variables at the $t$ and $\langle \cdot \rangle | Y(t) = X$ symbolizes a conditional averaging over the entire measurement period, where only measurements with $Y(t) = X$ are taken into account [2,7]. $D^{(1)}$ is the drift vector and $D^{(2)}$ the diffusion matrix.

Associated with the FPE (1) is a system of $N$ coupled Itô-Langevin equations, which can be written as [12,13]

$$\frac{dX}{dt} = h(X) + G(X)\Gamma(t) \tag{4}$$

Here $\Gamma(t)$ is a set of $N$ normally distributed random variables fulfilling

$$\langle \Gamma_i(t) \rangle = 0, \quad \langle \Gamma_i(t) \Gamma_j(t') \rangle = 2\delta_{ij}\delta(t-t') \tag{5}$$

that drives the stochastic evolution of $X$. The vectors $h$ and the matrices $G = \{g_{ij}\}$ for all $i, j = 1, \ldots, N$ are connected to the local drift and diffusion function through

$$D_i^{(1)}(X) = h_i(X) \text{ and } D_{ij}^{(2)}(X) = \sum_{k=1}^{N} g_{ik}(X)g_{jk}(X) \tag{6}$$

While the FPE describes the evolution of the joint distribution of the $N$ variables statistically, the system of Langevin equations in (4) models individual stochastic trajectories of the system. In Eq. (4) the term $h(X)$ contains the deterministic part of the macroscopic dynamics, while the locally linear functions $G(X)$ account for the amplitudes of the stochastic forces mirroring the different sources of fluctuations due to all sorts of microscopic interactions within the system. Here we should point out that we consider stationary processes, otherwise ensemble averages have to be taken [2,7].

The conditional moments in Eq. (3) are computed directly from data time series and are typically linear functions of $\Delta t$ for sufficiently small $\Delta t$ [2,7]. Then, through the limit in Eq. (2), both $D_i^{(1)}$ and $D_{ij}^{(2)}$ are determined.

The $N \times N$ matrix $G$ cannot be uniquely determined from the symmetric diffusion matrix $D^{(2)}$ for $N \geq 2$, because the number of unknown elements in $G$ exceeds the number of known elements in $D^{(2)}$ leading to $N^2 - \frac{1}{2}N(N + 1) = \frac{1}{2}N(N - 1)$ free parameters. However, a simple method to obtain $G$ from $D^{(2)}$ is the following. Due to its symmetry and positive semi-definiteness the diffusion matrix $D^{(2)}$ has only real, non-negative eigenvalues $\lambda_i (i = 1, \ldots, N)$. Therefore an orthogonal transformation $U$ can be found that diagonalizes $D^{(2)}$, i.e. $U^T D^{(2)} U = \text{diag}(\lambda_1, \ldots, \lambda_N)$. Taking the positive root of the eigenvalues and transforming it back we arrive at $g_{ij} = (\sqrt{D^{(2)}})_{ij}$ where the symbolic notation $\sqrt{D^{(2)}} = U \text{diag}(\sqrt{\lambda_1}, \ldots, \sqrt{\lambda_N})U^T$ is used for simplicity. General forms of $G$ can be constructed by multiplication of $\sqrt{D^{(2)}}$ with arbitrary orthogonal matrices. For our considerations the simple version is sufficient [12,7].

FIG. 1: Stream plot of the system in Eq. (8) with no stochasticity $(D^{(2)} = 0)$, i.e. $k_1 = k_2 = 0$, where the arrows formally represent the drift vector of synthetic data extracted from direct simulation of Eq. (8). The solid line represents a partial trajectory. Trajectories diverge from the unstable fixed point $r = 0$ and converge to the limit cycle $r = 1$. 

[Image 317x157 to 566x415]
In general the eigenvalues indicate the amplitude of the stochastic force and the corresponding eigenvector indicates the direction toward which such force acts. Even more interesting features, however, can be extracted from the eigenvalues and eigenvectors.

To each eigenvector of the diffusion matrix we can associate one independent source of stochastic forcing \( \Gamma_i \). In this scope, the eigenvectors can be regarded as defining principal axes for stochastic dynamics. For instance, the vector field aligned at each mesh point to the eigenvector associated to the smallest eigenvalue of matrix \( \mathbf{G} \) defines the paths in phase space towards which the fluctuations are minimal. Further, more, if the corresponding eigenvalues are very small compared to all the other ones at the respective mesh points, the corresponding stochastic forces can be neglected and the system has only \( N - 1 \) independent stochastic forces. In this case the problem can be reduced in one stochastic variable by an appropriate transformation of variables, since the eigenvectors in one coordinate system are the same as in another one (see Appendix A).

These are the central ideas of our study which we next apply to an analytical example, namely the Hopf-bifurcation.

### III. The Hopf-Bifurcation System with Stochastic Forcing: An Analytical Example

In what follows we consider the dynamical system

\[
\begin{align*}
\frac{dr}{dt} &= r(1 - r^2) + k_1 \Gamma_1 \\
\frac{d\theta}{dt} &= \alpha - r^2 + k_2 \Gamma_2
\end{align*}
\]  

Equation (8b) uses \( \alpha - r^2 \) instead of the usual \( 1 + r^2 \) in order to be able to avoid the systematic increase of \( \theta \) in time. For this reason we chose \( \alpha = \langle r^2 \rangle \).

Introducing non-zero stochastic terms \( k_1, k_2 \neq 0 \) for independent stochastic forces, \( \Gamma_1 \) and \( \Gamma_2 \), one has diagonal diffusion matrices with two eigenvalues, \( (k_1 r)^2 \) and \( k_3^2 \) at each point \( (r, \theta) \) corresponding to the radial and angular eigendirections respectively.

Since the eigenvectors do not change when changing the coordinate system (see Appendix A), “mixing” the radial and angular coordinates by e.g. changing to Cartesian should yield the same principal axis. We, therefore, now repeat our analysis for transformed variables \( x = r \cos \theta \) and \( y = r \sin \theta \).

Using Ito’s formulation for the transformation from polar to Cartesian coordinates, the Langevin system can be described by:

\[ \mathbf{X} = \mathbf{A} \mathbf{X} + \mathbf{B} \mathbf{w} \]
FIG. 4: Drift and diffusion functions for the system in Eq. (8) for the same conditions as in Fig. 2 obtained from analysis of the time series $x = r \cos \theta$ and $y = r \sin \theta$. For each drift and diffusion function the bottom surface reflects the numerical results, while the corresponding fitted surface on top is shifted upwards for clarity: (a) $\tilde{D}_{11}^{(2)}(x, y)$, (b) $\tilde{D}_{12}^{(2)}(x, y)$, (c) $\tilde{D}_{22}^{(2)}(x, y)$, (d) $\tilde{D}_{11}^{(1)}(x, y)$, (e) $\tilde{D}_{22}^{(1)}(x, y)$. The bar plots compare the coefficients that fit each function (dark bars) with the analytic, true ones (white bars) (see Tab. I).

defined through the drift and diffusion coefficients in transformed coordinates, $\tilde{h}_{ij}(x, y), \tilde{g}_{ij}(x, y)$ (see Appendix A):

$$\tilde{h}_{i} = \left( h_{i}/r - g_{22}^{2} \right) Y_{i} + (-1)^{i} h_{2} Y_{1} \delta_{i 1} + Y_{2} \delta_{i 1} \tag{9a}$$

$$\tilde{g}_{ij} = \left( \frac{g_{ij}}{r} (Y_{1} \delta_{i 1} + Y_{2} \delta_{i 2}) - g_{2j} (Y_{1} \delta_{j 1} + Y_{2} \delta_{j 1}) \right) \tag{9b}$$

with $i, j = 1, 2$ and $(Y_{1}, Y_{2}) = (x, y)$. Thus, from $\mathbf{G}$ in the polar coordinate system, Eq. (8) yields $\tilde{\mathbf{G}}$ for which the diffusion matrix $\tilde{\mathbf{D}}^{(2)} = \tilde{\mathbf{G}} \tilde{\mathbf{G}}^{T}$ reads

$$\tilde{\mathbf{D}}^{(2)} = \begin{bmatrix}
k_{1}^{2} y^{2} + k_{2}^{2} x^{2} & -k_{1}^{2} x y & k_{1}^{2} x y \\
-k_{2}^{2} x y & k_{2}^{2} x^{2} + k_{1}^{2} y^{2}
\end{bmatrix} \tag{10}$$

The eigenvalues and corresponding eigenvectors are $\tilde{\lambda}_{1} = k_{1}^{2} (x^{2} + y^{2})$ with $\tilde{\mathbf{v}}_{1} = (x/\sqrt{x^{2} + y^{2}}, y/\sqrt{x^{2} + y^{2}})$ and $\tilde{\lambda}_{2} = k_{2}^{2} (x^{2} + y^{2})$ with $\tilde{\mathbf{v}}_{2} = (-y/\sqrt{x^{2} + y^{2}}, x/\sqrt{x^{2} + y^{2}})$.

In contrast to the eigendirections, the eigenvalues depend on the Jacobian of our transformation, since the nonlinear
transformation from polar to Cartesian coordinates changes the metric. According to Appendix A the eigenvalues in polar coordinates are
\[ \lambda_i = \lambda_i/s_i^2 \]
for \( i = 1, 2 \) with \( s_1^2 = (\partial_r^2 + \partial_{\theta}^2) = 1 \) and \( s_2^2 = (\partial_r^2 + \partial_{\theta}^2) = r^2 \). The eigenvalues of our system in Cartesian coordinates are shown as solid and dashed lines in Fig. 2 and a full derivation for the eigenvalues in different coordinate systems is given in Appendix A for the general \( N \)-dimensional case.

Notice that, if we determine the eigenvalues in a Cartesian system, the metric is Euclidean, i.e. the eigenvalues measured in the same units in both \( x \)- and \( y \)-direction directly characterize the diffusion in principal directions. Nonlinear transformations of coordinates such as the one from Cartesian to polar coordinates, however, generally change the metric. In such a system the direction of the maximal eigenvalue is not necessarily the direction with the highest diffusion. For example, in the Hopf-bifurcation in Eqs. (8), the maximal eigenvalue for \( r < k_2/k_1 \) is in azimuthal direction, but the maximal diffusion is still in the radial direction.

Looking to \( x \) and \( y \) as measured variables, they define our Cartesian system and therefore \( \lambda_1 \) is our maximal eigenvalue and \( \lambda_2 \) is associated with the eigendirection showing minimal stochastic fluctuation.

We proceed to verify this result by numerically analyzing the time series \( x \) and \( y \) plotted in Fig. 3 following the procedure described in the previous section. The procedure is as follows. First, we compute the time series according to Eq. (8) in coordinates \( r \) and \( \theta \). Then, we transfer this time series to Cartesian coordinates \( x, y \). In the new coordinates we calculate the drift vectors \( \tilde{D}^{(1)} \) and diffusion matrices \( \tilde{D}^{(2)} \). Finally, having matrices \( \tilde{D}^{(2)} \) at each mesh point, one easily computes the eigenvalues and eigenvectors, as shown in Fig. 3. The analytical results given by Eq. (10) are reproduced nicely: the stochastic contribution has two eigendirections, one in the radial direction and another in the angular one; the larger eigenvalue is associated with the radial direction. The deviations observed for the minimum eigenvalue occur due to differences in the quality of statistics in different regions of phase space (value of \( r \)): in particular regions in phase space at intermediate values of \( r \) are less frequently realized in the data series than others, resulting in a reduced accuracy of the estimated eigenvalues in these regions.

Figures 4a-c show the components of the diffusion matrix, namely \( \tilde{D}^{(1)}_1, \tilde{D}^{(2)}_1 \) and \( \tilde{D}^{(2)}_2 \), as a function of both \( x \) and \( y \). Notice that the diffusion matrix is symmetric (see Eq. (3)) and therefore \( \tilde{D}^{(2)}_1 = \tilde{D}^{(2)}_2 \). The surfaces obtained numerically are fitted to a full quadratic polynomial in \( x \) and \( y \) through a least square procedure. The coefficients of these polynomial fits are also shown in Fig. 4a-c. The results of Fig. 4a-c comply with the analytical expression in Eq. (10). A similar analysis is done for the drift vector, i.e. for the function \( \tilde{D}^{(1)}_1 \) and \( \tilde{D}^{(1)}_2 \), shown in Figs. 4d-e. Table I lists the values obtained for the coefficients of the surface fits in Fig. 4.

It is worth stressing that the aforementioned procedure can likewise be applied to measured multi-dimensional time series in cases where no additional information on the underlying dynamics is available. In this case the corresponding fields of eigenvectors indicate the directions with largest and smallest stochastic contributions. Using this information, a transform of variables to a coordinate system aligned with the direction of the smallest stochastic contribution can be applied, which in the present case would be the tangential direction. In the representation of these new coordinates, the stochastic contribution will then be reduced.

In the previous case, the two eigenvectors for diffusion were introduced tangentially and perpendicularly to the limit cycle. If the principal axes for the stochastic contribution are not aligned with the trajectories of the deterministic part of the system, one might suspect an interplay between both stochastic contribution and the drift of the system.

In order to demonstrate the wider applicability of our method, we next consider such a case, with the major axis aligned with \( y \)-direction. We therefore investigate the system

\[
\frac{dr}{dt} = r(1 - r^2) + K_1 r \cos \theta \Gamma_1 + K_2 r \sin \theta \Gamma_2
\]

\[
\frac{d\theta}{dt} = (\alpha - r^2) - K_1 \sin \theta \Gamma_1 + K_2 \cos \theta \Gamma_2
\].

In Cartesian coordinates the dynamics are described by

\[
\frac{dx}{dt} = \tilde{h}_1 + K_1 r \Gamma_1
\]

\[
\frac{dy}{dt} = \tilde{h}_2 + K_2 r \Gamma_2.
\]

with

\[
\tilde{h}_1 = x(1 - (x^2 + y^2)) - y(\alpha - (x^2 + y^2)) + \frac{x}{x^2 + y^2}((K_1^2 - 2K_2^2)y^2 - K_2^2 x^2)
\]

\[
\tilde{h}_2 = y(1 - (x^2 + y^2)) - x(\alpha - (x^2 + y^2))
\]

TABLE I: Coefficients of quadratic forms describing the surfaces in Fig. 2

| Coefficients for | 1 | x | x^2 | x^3 | y | y^2 | y^3 | xy | x^2 y | xy^2 |
|------------------|---|---|-----|-----|---|-----|-----|-----|-------|-------|
| \( D^{(1)}_1 \) | -0.009 | 1.00 | 0.077 | -1.01 | 0.820 | -0.012 | 1.26 | 0.014 | 1.23 | 1.24 |
| \( D^{(1)}_2 \) | -0.034 | 0.761 | 0.133 | -1.09 | 0.99 | 0.036 | -1.22 | 0.101 | -1.00 | -1.21 |
| \( D^{(2)}_1 \) | 0.005 | -0.008 | 0.262 | \( \sim 0 \) | 0.001 | 0.013 | \( \sim 0 \) | \( \sim 0 \) | \( \sim 0 \) | \( \sim 0 \) |
| \( D^{(2)}_2 \) | \( \sim 0 \) | -0.003 | 0.14 | \( \sim 0 \) | -0.004 | -0.011 | \( \sim 0 \) | 0.255 | \( \sim 0 \) | \( \sim 0 \) |
| \( D^{(2)}_3 \) | 0.005 | \( \sim 0 \) | -0.004 | \( \sim 0 \) | -0.006 | 0.255 | \( \sim 0 \) | \( \sim 0 \) | \( \sim 0 \) | \( \sim 0 \) |
FIG. 5: Hopf-system with strong stochastic contribution along the $y$-direction. Drift functions are (a) $\tilde{D}_1^{(1)}(x, y)$ and (b) $\tilde{D}_2^{(1)}(x, y)$. (c) Probability density function of $(x, y)$ shows the most visited regions in phase space. Diffusion functions are (d) $\tilde{D}_{11}^{(2)}(x, y)$, (e) $\tilde{D}_{12}^{(2)}(x, y) = D_{21}^{(2)}(x, y)$ and (f) $\tilde{D}_{22}^{(2)}(x, y)$. Results were obtained by analyzing the time series $x = r \cos \theta$ and $y = r \sin \theta$ according to Eq. (11) with the same upward shifting for clarity as in Fig. 4. The probability density function in (c) explains the deviations observed for the diffusion functions (see text).

\[ + \frac{y}{x^2 + y^2} \left( (K_2^2 - 2K_1^2)x^2 - k_1^2y^2 \right). \] (13b)

Here we chose $K_1 = 0.05$ and $K_2 = 0.5$. In other words, the impact of stochastic fluctuations is large in the $y$-direction and small along the $x$-axis. Notice that the determinant of the diffusion matrix is not preserved, since the Jacobian of the transformation is not the identity matrix (see Append. A).

As already done in the previous example, we integrate system (11) numerically and analyze the resulting data series of $x$ and $y$ values. As can be seen from Fig. 5, the drift functions $\tilde{h}_1$ and $\tilde{h}_2$ are properly derived (see Figs. 5a and 5b), as well as the terms of the diffusion matrices $\tilde{D}^{(2)}$ (Figs. 5d-f). The deviations observed for the small $\tilde{D}_{11}^{(2)}$ at particular regions of phase space (Fig. 5b) are due to the lack of accurate statistics at those regions (see the PDF in Fig. 5c).

Interestingly, in both cases presented above, our analysis estimates better the diffusion matrix than the corresponding drift vector, contrary to what is known in other situations [2]. The reason for the poor quality of the drift estimate is probably due to the small time increment used for the estimate, which spoils particularly the result for the drift. See Ref. [19] for details.

As can be seen in Fig. 6a and 6b, the eigendirections are properly derived as well as the two eigenvalues plotted in Fig. 6c and 6d as a function of $r$ and $\theta$, respectively.

With these two examples we first have illustrated that our numerical approach succeeds in estimating the drift and the diffusion functions contributions in a two-dimensional stochastic system. A straightforward computation of the eigenvectors and eigenvalues of the diffusion matrices then enables us to ascertain a set of principal stochastic directions. We would like to note that the method described here bears some resemblance to the well-known principal component analysis [20]. In contrast to principal component analysis performed on the distribution of the measured data the method we propose, however, focuses on the stochastic fluctuations in time. The properties of these fluctuations in general are dependent on the position $X$ resulting in vector fields of the principal axes of the stochastic contributions to the dynamics of the system of consideration. The two methods seem however related and should provide complementary insights when applied to specific sets of empirical data. Such matters are beyond the scope of this paper and will be addressed elsewhere.

We also emphasize that the minimal eigenvalue of the diffusion matrix not necessarily points towards the direction the systems as a whole can be expected to evolve to. It rather reflects the direction of minimal stochastic forcing only. For prediction of the system’s evolution in general time propagators need to be taken into account which involve both the stochastic (diffusion) and the deterministic (drift) parts of the dynamics [2]. Nevertheless, being able to ascertain the directions in phase space where fluctuations are weaker enables one
FIG. 6: Results of the analysis of diffusion matrices plotted in Fig. 5, obtained from a system with an increased level of dynamical noise on the y component (second example). The individual panels depict (a) the vector fields of eigenvectors for the minimum and (b) the maximum eigenvalue, normalized to the respective eigenvalues, and the dependence of the eigenvalues both on (c) r and (d) θ.

to choose a transformation such that part of the new variables have small stochastic terms, reducing the number of variables which are affected by stochastic forces. Further, this is also of physical interest to see in which variables noise is acting, for example to which variable a thermal bath is physically connected, and it may be of technical interest to better detect a noise source, for example in an electric circuit [21].

IV. DISCUSSION AND CONCLUSIONS

In this paper we introduce the concept of eigendirections for the stochastic dynamics in systems of arbitrary dimension. The procedure builds on the modeling of complex systems by means of drift and diffusion functions, that specify a system of coupled Langevin equations. Estimates for these functions can be obtained directly from measurements on the systems without prior knowledge on its dynamics following an approach described in Ref. [1, 2].

In Langevin systems the stochastic forcing is composed of independent Gaussian δ-correlated stochastic fluctuations. The way how these fluctuations effect the system depends both on the diffusion matrix and on the coordinate system chosen.

Whereas in former publications much attention has been paid to imperfections of the Langevin process, like measurement noise or correlated noise and finite size sampling [22–25], here we introduced a method which allows us to determine the eigendirections along which each stochastic force acts. Each direction of forcing is defined through the eigenvector and the corresponding eigenvalue accounting for its amplitude. The set of eigenvectors does not depend on the choice of the coordinate system and is therefore characteristic for the system. Further, for the particular case where a number k of eigenval-
uses are negligible in comparison with the others at each mesh point, the number of stochastic variables can be reduced. Even in cases where the number of stochastic variables cannot be reduced, the eigenvector associated with the lowest eigenvalues at each point in phase space indicates the path with minimal stochastic forcing. In any case a transform to new coordinates in the directions of minimum stochastic forcing can be performed, increasing the relative amplitude of the deterministic components and the predictability of the corresponding variable.

The method was successfully applied to a two-dimensional system exhibiting a Hopf-bifurcation. For the Hopf-bifurcation the diffusion matrix is better estimated through our analysis than the corresponding drift vector, contrary to what is known in many other situations [3].

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Appendix A: The diffusion matrix in an arbitrary coordinate system

Consider a transformation of variables \( X = \{X_i\} \rightarrow \tilde{X} = \{\tilde{X}_i\} \) with \( i = 1, \ldots, N \), which is given by a two-times continuously differentiable deterministic vector function \( \tilde{F} \)

\[
\tilde{X} = \tilde{F}(X, t) .
\]  
(A1)

Using Itô’s formula [12][13], a truncated form of the Itô Taylor expansion, the Langevin equations for the new variables take the form

\[
d\tilde{X}_i = \frac{\partial F_i}{\partial t} dt + \sum_{k=1}^{N} \frac{\partial F_i}{\partial X_k} dX_k + \frac{1}{2} \sum_{k=1}^{N} \sum_{l=1}^{N} \frac{\partial^2 F_i}{\partial X_k \partial X_l} d\Gamma_{kl} + \ldots .
\]  
(A2)

Eq. (A1) can be rewritten to

\[
dX_i = D_i^{(1)}(X) dt + \sum_{j=1}^{N} g_{ij}(X) \Gamma_j(t) \sqrt{dt} .
\]  
(A3)

Inserting this expression into Eq. (A2), retaining all terms in the expansion up to order \( dt \) while neglecting terms of order \((dt)^{3/2}\) or higher, and taking advantage of the statistics of the stochastic forcing, Eq. (5), one obtains for the third expression on the r.h.s.

\[
\frac{1}{2} \sum_{k=1}^{N} \sum_{l=1}^{N} \left[ (D_k^{(1)})_{kl} + \sum_{j=1}^{N} g_{kj} \Gamma_j(t) \sqrt{dt} \right] \frac{\partial^2 F_i}{\partial X_k \partial X_l} dt + \sum_{m=1}^{N} g_{im} \Gamma_m(t) \sqrt{dt} = \frac{1}{2} \sum_{k=1}^{N} \sum_{l=1}^{N} \sum_{j,m=1}^{N} 2g_{kj}g_{lm} \frac{\partial^2 F_i}{\partial X_k \partial X_l} \Gamma_{jm} dt .
\]  
(A4)

Thus, the transformed Langevin equation has the form

\[
d\tilde{X}_i = \frac{\partial \tilde{F}_i}{\partial t} dt + \sum_{k=1}^{N} \frac{\partial \tilde{F}_i}{\partial X_k} (D_k^{(1)})_{kl} dt + \sum_{j=1}^{N} g_{ij} \Gamma_j(t) \sqrt{dt} + \sum_{k=1}^{N} \sum_{l=1}^{N} \sum_{j=1}^{N} g_{lj} \frac{\partial^2 \tilde{F}_i}{\partial X_k \partial X_l} dt + \mathcal{O}(dt)^{3/2} .
\]  
(A5)

Restricting to stationary processes (\( \frac{\partial \tilde{F}_i}{\partial t} = 0 \)) and using the notation

\[
\frac{d\tilde{X}_i}{dt} = \tilde{D}_i^{(1)}(\tilde{X}) + \sum_{j=1}^{N} \tilde{g}_{ij}(\tilde{X}) \Gamma_j(t) .
\]  
(A6)

the transformed drift function reads

\[
\tilde{D}_i^{(1)} = \sum_{k=1}^{N} \left[ (D_k^{(1)})_{ik} + \sum_{l=1}^{N} g_{lj} \frac{\partial^2 \tilde{F}_i}{\partial \tilde{X}_k \partial \tilde{X}_l} \right] .
\]  
(A7)

The components of matrix \( G \) transform as

\[
\tilde{g}_{ij} = \sum_{k=1}^{N} g_{kj} \frac{\partial \tilde{F}_i}{\partial \tilde{X}_k} \quad \Rightarrow \quad \tilde{G} = JG ,
\]  
(A8)

where \( J \) is the Jacobian of our transformation from coordinates \( X_k \) (vector basis \( e_j \)) to \( \tilde{X}_i \) (vector basis \( \tilde{e}_j \)).

From elementary algebra it is known that a vector \( v \) can be written in both coordinate systems:

\[
v = \sum_{j} \tilde{v}_j \tilde{e}_j = \sum_{j} v_j \sum_{k} J_{jk} e_k
\]
Thus, the matrix \( \tilde{\textbf{U}} \) incorporating the columns of eigenvectors \( \textbf{u}_k \) of matrix \( \tilde{\textbf{G}} \), with coordinates in basis \( \tilde{e}_i \), can be written as \( \textbf{U} = \textbf{U}^T \tilde{\textbf{U}} \) with \( \textbf{U} = [\textbf{u}_1 \ \textbf{u}_2 \ldots \ \textbf{u}_N] \) fulfilling \( \textbf{U}^{-1} = \textbf{U}^T \). Thus, from Eq. (A8) one obtains

\[
\tilde{\textbf{U}}^T \tilde{\textbf{G}} \tilde{\textbf{G}}^T \tilde{\textbf{U}} = \textbf{U}^T \textbf{G} \textbf{G}^T \textbf{U}.
\]  

(A10)

So, the eigenvectors of \( \textbf{D}^{(2)} \) and \( \tilde{\textbf{D}}^{(2)} \) are the same, apart from the basis in which they are considered, and in the original (Cartesian) coordinate system \( \textbf{X} \), the transformation to principal axes is given by Eq. (A10) i.e.

\[
\textbf{U}^T \textbf{D}^{(2)} \textbf{U} = \text{diag}[\lambda_1 \ \lambda_2 \ldots \ \lambda_N] \quad (A11)
\]

At regular points the transformation in Eq. (A1) can be inverted,

\[
\textbf{X} = \textbf{F}^{-1}(\tilde{\textbf{X}}) =: \textbf{f}(\tilde{\textbf{X}}). \quad (A12)
\]

By definition \( \textbf{f}(\tilde{\textbf{X}}) \) is chosen such that the normalized eigenvectors are given by

\[
\textbf{u}_k = \frac{1}{s_k} \left. \frac{\partial \textbf{f}}{\partial X_k} \right| \quad \text{with} \quad s_k = \left| \frac{\partial \textbf{f}}{\partial X_k} \right|. \quad (A13)
\]

The Jacobian of \( \textbf{f}(\tilde{\textbf{X}}) \) can then be written as

\[
\tilde{\textbf{J}}_{ik}(\tilde{\textbf{X}}) = \frac{\partial f_i}{\partial X_k}, \quad (A14)
\]

or in a more convenient form

\[
\tilde{\textbf{J}} = [\textbf{u}_1 \ \textbf{u}_2 \ldots \ \textbf{u}_N] \text{diag}[s_1 \ s_2 \ldots \ s_N] =: \textbf{Us}. \quad (A15)
\]

The diagonal matrix \( \textbf{s} \) describes the metric of the new system. The Jacobian of the transformation in Eq. (A1) is the inverse relation

\[
\textbf{J} = \textbf{s}^{-1} \textbf{U}^T. \quad (A16)
\]

Introducing this relation into Eq. (A8) yields

\[
\tilde{\textbf{G}} = \textbf{JG} = \textbf{s}^{-1} \textbf{U}^T \textbf{G}, \quad (A17)
\]

i.e. the transformed diffusion matrix reads

\[
\tilde{\textbf{D}}^{(2)} = \tilde{\textbf{G}} \tilde{\textbf{G}}^T = \textbf{s}^{-1} \textbf{U}^T \textbf{G} \textbf{G}^T \textbf{Us}^{-1}. \quad (A18)
\]

Inserting Eqs. (A11) and (A15) finally leads to the diagonal matrix

\[
\tilde{\textbf{D}}^{(2)} = \text{diag} \left[ \frac{\lambda_1}{s_1} \frac{\lambda_2}{s_2} \ldots \frac{\lambda_N}{s_N} \right]. \quad (A19)
\]

FIG. 7: Interchanging between coordinate systems. Arrows indicate that eigenvectors of the eigenvalues of the diffusion matrix (see Eq. (A19)) at each point of the mesh from which drift and diffusion coefficients are derived. Here the two-dimensional case is illustrated. Extension to arbitrary number of variables is straightforward.

In practice, we have no access to the transformation in Eq. (A1). Instead we have a grid (numerical) representation in \( \textbf{X} \) coordinate system and search for the transformation in Eq. (A1) is such that one of the new variables, say \( \tilde{X}_1 \), is defined by the field of eigenvectors \( \textbf{u}_1 \) corresponding to the maximal eigenvalue \( \lambda_1 \). Then, each orthogonal direction is consequently defined by one of the other (orthogonal) eigenvectors \( \textbf{u}_j \) with \( j = 1, \ldots, N \). Figure 7 illustrates this for the two-dimensional case.

To obtain the coordinates of each point \( P \) in this grid in the new coordinates \( \tilde{\textbf{X}} \) one considers the transformation (A1) together with the inverted transformation (A12) and solves the system of PDEs in (A13). The problem lies in the coupling introduced by the factors \( s_k \) and the additional complication of finding the correct boundary conditions which turns the solution of the PDEs into a complicated problem, even numerically, although in 2D the \( s_k \) can be neglected as we are only interested in finding the direction \( \frac{\partial f_i}{\partial X_k} \). In most cases, however, it should be possible to guess a suitable transform from visual inspection of the fields of eigenvectors. From one point to the next one in the mesh, the eigenvectors are sorted according to continuity arguments. Further, depending on the obtained fields of eigenvectors, a scaled polar form or hyperbolic coordinates may be considered, where the parameters of said transform can then be fitted to the vector fields. Such particular cases depend on the specific data set at hand and will be addressed elsewhere.
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[27] In practice binning or kernel based approaches with a certain threshold are applied in order to evaluate the condition $Y(t) = X$. See e.g. Ref. [14] for details.