Learning Stochastic Dynamical Systems via Bridge Sampling

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Abstract

We develop algorithms to automate discovery of stochastic dynamical system models from noisy, vector-valued time series. By discovery, we mean learning both a nonlinear drift vector field and a diagonal diffusion matrix for an Itô stochastic differential equation in $\mathbb{R}^d$. We parameterize the vector field using tensor products of Hermite polynomials, enabling the model to capture highly nonlinear and/or coupled dynamics. We solve the resulting estimation problem using expectation maximization (EM). This involves two steps. We augment the data via diffusion bridge sampling, with the goal of producing time series observed at a higher frequency than the original data. With this augmented data, the resulting expected log likelihood maximization problem reduces to a least squares problem. Through experiments on systems with dimensions one through eight, we show that this EM approach enables accurate estimation for multiple time series with possibly irregular observation times. We study how the EM method performs as a function of the noise level in the data, the volume of data, and the amount of data augmentation performed.

Traditional mathematical modeling in the sciences and engineering often has as its goal the development of equations of motion that describe observed phenomena. Classically, these equations of motion usually took the form of deterministic systems of ordinary or partial differential equations (ODE or PDE, respectively). Especially in systems of contemporary interest in biology and finance where intrinsic noise must be modeled, we find stochastic differential equations (SDE) used instead of deterministic ones. Still, these models are often built from first principles, after which the model's predictions (obtained, for instance, by numerical simulation) are compared against observed data.

Recent years have seen a surge of interest in using data to automate discovery of ODE, PDE, and SDE models. These machine learning approaches complement traditional modeling efforts, using available data to constrain the space of plausible models, and shortening the feedback loop linking model development to prediction and comparison to real observations. We posit two additional reasons to develop algorithms to learn SDE models. First, SDE models—including the models considered here—have the capacity to model highly nonlinear, coupled stochastic systems, including systems whose equilibria are non-Gaussian and/or multimodal. Second, SDE models often allow for interpretability. Especially if the terms on the right-hand side of the SDE are expressed in terms of commonly used functions (such as polynomials), we can obtain a qualitative understanding of how the system's variables influence, regulate, and/or mediate one other.

In this paper, we develop an algorithm to learn SDE models from high-dimensional time series. To our knowledge, this is the most general expectation maximization (EM) approach to learning an SDE with multidimensional drift vector field and diagonal diffusion matrix. Prior EM approaches were restricted to one-dimensional SDE [**?**], or used a Gaussian process approximation, linear drift approximation, and approximate maximization [**?**]. To develop our method, we use diffusion bridge sampling as in [6, 17], which focused on Bayesian nonparametric methods for SDE in $\mathbb{R}^1$. After

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augmenting the data using bridge sampling, we are left with a least-squares problem, generalizing the work of [3] from the ODE to the SDE context.

In the literature, variational Bayesian methods are the only other SDE learning methods that have been tested on high-dimensional problems [**?**]. These methods use approximations consisting of linear SDE with time-varying coefficients [**?**], kernel density estimates [**?**], or Gaussian processes [**?**]. In contrast, we parameterize the drift vector field using tensor products of Hermite polynomials; as mentioned above, the resulting SDE has much higher capacity than linear and/or Gaussian process models.

Many other techniques explored in the statistical literature focus on scalar SDE [1, 7, 8, 18].

As mentioned, differential equation discovery problems have attracted considerable recent interest. A variety of methods have been developed to learn ODE [3, 13, 15**? ? ?** , 16] as well as PDE [12, 14**? ?**]. Unlike many of these works, we do not focus on model selection; if needed, our methods can be combined with model selection procedures developed in the ODE context [4, 5].

Points to cover:

- data specification, DONE
- Hermite polynomial and drift function representation, DONE
- Expectation and maximization formulas assuming data is filled in, CLOSE2DONE
- Filling data in with diffusion bridge, DONE
- MCMC iterations of brownian bridge using girsanov likelihood, DONE
- how synthetic data is generated
- results: 1D, 2D, 3D damped duffing, 3D lorenz
- plots: error of theta vs noise, error vs amount of data (number of data points) parametric curves for noise levels, brownian bridge plots for illustration, ...

1 Problem Setup

Let W_t denote Brownian motion in $\mathbb{R}^d$—informally, an increment dW_t of this process has a multivariate normal distribution with zero mean vector and covariance matrix $I dt$. Let X_t denote an $\mathbb{R}^d$-valued stochastic process that evolves according to the Itô SDE

$$dX_t = f(X_t)dt + \Gamma dW_t. \tag{1}$$

For rigorous definitions of Brownian motion and SDE, see [2, 19]. The nonlinear vector field $f : \Omega \subset \mathbb{R}^d \to \mathbb{R}^d$ is the *drift* function, and the $d \times d$ matrix Γ is the *diffusion* matrix. To reduce the number of model parameters, we assume $\Gamma = \operatorname{diag}\gamma$.

Our goal is to develop an algorithm that accurately estimates the functional form of f and the vector γ from time series data.

Parameterization. We parameterize f using Hermite polynomials. The n-th Hermite polynomial takes the form

$$H_n(x) = (\sqrt{2\pi}n!)^{-1/2}(-1)^n e^{x^2/2}\frac{\mathrm{d}^n}{\mathrm{d}x^n}e^{-x^2/2} \tag{2}$$

Let $\langle f, g\rangle_w = \int_{\mathbb{R}} f(x)g(x)\exp(-x^2/2)\,dx$ denote a weighted L^2 inner product. Then, $\langle H_i, H_j\rangle_w = \delta_{ij}$, i.e., the Hermite polynomials are orthonormal with respect to the weighted inner product. In fact, with respect to this inner product, the Hermite polynomials form an orthonormal basis of $L^2_w(\mathbb{R}) = \{f \,:\, \langle f, f\rangle_w < \infty\}$.

Now let $\alpha = (\alpha_1, \ldots, \alpha_d) \in \mathbb{Z}^d_+$ denote a multi-index. We use the notation $|\alpha| = \sum_j \alpha_j$ and $x^\alpha = \prod_j (x_j)^{\alpha_j}$ for $x = (x_1, \ldots, x_d) \in \mathbb{R}^d$. For $x \in \mathbb{R}^d$ and a multi-index α, we also define

$$H_\alpha(x) = \prod_{j=1}^{d} H_{\alpha_j}(x_j). \tag{3}$$

We write $f(x) = (f_1(x), \dots f_d(x))$ and then parameterize each component:

$$f_j(x) = \sum_{m=0}^{M} \sum_{|\alpha|=m} \beta_\alpha^j H_\alpha(x). \tag{4}$$

We see that the maximum degree of $H_\alpha(x)$ is $|\alpha|$. Hence we think of the double sum in (4) as first summing over degrees and then summing over all terms with a fixed maximum degree. We say maximum degree because, for instance, $H_2(z) = (z^2 - 1)/(\sqrt{2\pi}2)^{1/2}$ contains both degree 2 and degree 0 terms.

There are $\binom{m+d-1}{d-1}$ possibilities for a d-dimensional multi-index α such that $|\alpha| = m$. Summing this from $m = 0$ to M, there are $\widetilde{M} = \binom{M+d}{d}$ total multi-indices in the double sum in (4). Let (i) denote the i-th multi-index according to some ordering. Then we can write

$$f_j(x) = \sum_{i=1}^{\widetilde{M}} \beta_{(i)}^j H_{(i)}(x). \tag{5}$$

Data. We consider our data $\mathbf{x} = \{x_j\}_{j=0}^L$ to be direct observations of X_t at discrete points in time $\mathbf{t} = \{t_j\}_{t=0}^L$. Note that these time points do not need to be equispaced.

To achieve our estimation goal, we apply expectation maximization (EM). We regard $\mathbf{x}$ as the incomplete data. Let $\Delta t = \max_j (t_j - t_{j-1})$ be the maximum interobservation spacing. We think of the missing data $\mathbf{z}$ as data collected at a time scale $h \ll \Delta t$ that is fine enough such that the transition density of (1) is approximately Gaussian. To see how this works, let $\mathcal{N}(\mu, \Sigma)$ denote a multivariate normal with mean vector μ and covariance matrix Σ. Now discretize (1) in time via the Euler-Maruyama method with time step $h > 0$; the result is

$$\widetilde{X}_{n+1} = \widetilde{X}_n + f(\widetilde{X}_n)h + h^{1/2}\Gamma Z_{n+1}, \tag{6}$$

where $Z_{n+1} \sim \mathcal{N}(0, I)$ is a standard multivariate normal, independent of X_n. This implies that

$$(\widetilde{X}_{n+1} | \widetilde{X}_n = v) \sim \mathcal{N}(v + f(v)h, h\Gamma^2). \tag{7}$$

As h decreases, $\widetilde{X}_{n+1} | \widetilde{X}_n = v$—a Gaussian approximation—will converge to the true transition density $X_{(n+1)h} | X_{nh} = v$, where X_t refers to the solution of (1).

Diffusion Bridge. To augment or complete the data, we employ diffusion bridge sampling, using a Markov chain Monte Carlo (MCMC) method that goes back to [10, 11]. Let us describe our version here. We suppose our current estimate of $\boldsymbol{\theta} = (\beta, \gamma)$ is given. Define the diffusion bridge process to be (1) conditioned on both the initial value x_i at time t_i, and the final value x_{i+1} at time t_{i+1}. The goal is to generate sample paths of this diffusion bridge. By a sample path, we mean $F - 1$ *new* samples $\{z_{i,j}\}_{j=1}^{F-1}$ at times $t_i + jh$ with $h = (t_{i+1} - t_i)/F$.

To generate such a path, we start by drawing a sample from a Brownian bridge with the same diffusion as (1). That is, we sample from the SDE

$$d\widehat{X}_t = \Gamma dW_t \tag{8}$$

conditioned on $\widehat{X}_{t_i} = x_i$ and $\widehat{X}_{t_{i+1}} = x_{i+1}$. This Brownian bridge can be described explicitly:

$$\widehat{X}_t = \Gamma(W_t - W_{t_i}) + x_i - \frac{t - t_i}{t_{i+1} - t_i}(\Gamma(W_{t_{i+1}} - W_{t_i}) + x_i - x_{i+1}) \tag{9}$$

Here $W_0 = 0$ (almost surely), and $W_t - W_s \sim \mathcal{N}(0, (t - s)I)$ for $t > s \geq 0$.

Let $\mathbb{P}$ denote the law of the diffusion bridge process, and let $\mathbb{Q}$ denote the law of the Brownian bridge (9). Using Girsanov's theorem [9], we can show that

$$\frac{d\mathbb{P}}{d\mathbb{Q}} = C \exp\left(\int_{t_i}^{t_{i+1}} f(\widehat{X}_s)^T \Gamma^{-2} \, d\widehat{X}_s - \frac{1}{2} \int_{t_i}^{t_{i+1}} f(\widehat{X_s})^T \Gamma^{-2} f(\widehat{X_s}) \, ds \right), \tag{10}$$

where the constant C depends only on x_i and x_{i+1}. The left-hand side is a Radon-Nikodym derivative, equivalent to a density or likelihood; the ratio of two such likelihoods is the accept/reject ratio in the Metropolis algorithm [Stuart 2010].

Putting the above pieces together yields the following Metropolis algorithm to generate diffusion bridge sample paths. Fix $F \geq 2$ and $i \in \{0, \dots, L-1\}$. Assume we have stored the previous Metropolis step, i.e., a path $\mathbf{z}^{(\ell)} = \{z_{i,j}\}_{j=1}^{F-1}$.

1. Use (9) to generate samples of $\widehat{X}_t$ at times $t_i + jh$, for $j = 1, 2, \dots, F-1$ and $h = (t_{i+1} - t_i)/F$. This is the proposal $\mathbf{z}^* = \{z^*_{i,j}\}_{j=1}^{F-1}$.
2. Numerically approximate the integrals in (10) to compute the likelihood of the proposal. Specifically, we compute
$$p(\mathbf{z}^*)/C = \sum_{j=0}^{F-1} f(z^*_{i,j})^T \Gamma^{-2} (z^*_{i,j+1} - z^*_{i,j}) - \frac{h}{4} \sum_{j=0}^{F-1} \left[f(z^*_{i,j})^T \Gamma^{-2} f(z^*_{i,j}) + f(z^*_{i,j+1})^T \Gamma^{-2} f(z^*_{i,j+1}) \right]$$
3. Accept the proposal with probability $p(\mathbf{z}^*)/p(\mathbf{z}^{(\ell)})$—note the factors of C cancel. If the proposal is accepted, then set $\mathbf{z}^{(\ell+1)} = \mathbf{z}^*$. Else set $\mathbf{z}^{(\ell+1)} = \mathbf{z}^{(\ell)}$.

We initialize this algorithm with a Brownian bridge path, run for 10 burn-in steps, and then use subsequent steps as the diffusion bridge samples we seek.

Expectation Maximization (EM). Let us now give details to justify the intuition expressed above, that employing the diffusion bridge to augment the data on a fine scale will enable estimation. Let $\mathbf{z}$ a diffusion bridge sample path. We interleave this sampled data together with the observed data $\mathbf{x}$ to create the completed time series
$$\mathbf{y}^{(r)} = \{y_j^{(r)}\}_{j=1}^{N},$$
where $N = LF + 1$. By interleaving, we mean that $y^{(r)}_{1+iF} = x_i$ for $i = 0, 1, \dots, L$, and that $y^{(r)}_{1+j+iF} = z_{i,j}$ for $j = 1, 2, \dots, F-1$ and $i = 0, 1, \dots, L-1$.

Let h_j denote the elapsed time between observations y_j and y_{j+1}. Using the completed data, the temporal discretization (6) of the SDE, the Markov property, and property (7), we have:
$$\begin{aligned}
\log p(\mathbf{x}, \mathbf{z} \mid \boldsymbol{\theta}) &= \log p(\mathbf{y} \mid \boldsymbol{\theta}) \\
&= \sum_{j=1}^{N-1} \log p(y_{j+1} \mid y_j, \boldsymbol{\theta}) \\
&= -\sum_{j=1}^{N-1} \left[\sum_{i=1}^{d} \frac{1}{2} \log(2\pi h_j \gamma_i^2) \right. \\
&\quad \left. + \frac{1}{2h_j} (y_{j+1} - y_j - h_j \sum_{k=1}^{\widetilde{M}} \beta_{(k)} H_{(k)}(y_j))^T \Gamma^{-2} (y_{j+1} - y_j - h_j \sum_{\ell=1}^{\widetilde{M}} \beta_{(\ell)} H_{(\ell)}(y_j)) \right].
\end{aligned}$$

EM. The EM algorithm consists of two steps, computing the expectation of the log likelihood function (on the completed data) and then maximizing it with respect to the parameters

1. Start with an initial guess for the parameters, $\boldsymbol{\theta}^{(0)}$.
2. For the expectation (or E) step,
$$Q(\boldsymbol{\theta}, \boldsymbol{\theta}^{(k)}) = \mathbb{E}_{\mathbf{z}|\mathbf{x},\boldsymbol{\theta}^{(k)}} [\log p(\mathbf{x}, \mathbf{z} \mid \boldsymbol{\theta})] \tag{11}$$
Our plan is to evaluate this expectation via bridge sampling. That is, we will sample from diffusion bridges $\mathbf{z} \mid \mathbf{x}, \boldsymbol{\theta}^{(k)}$. Then $(\mathbf{x}, \mathbf{z})$ will be a combination of the original data together with sample paths.

3. For the maximization (or M) step, we start with the current iterate and a dummy variable $\boldsymbol{\theta}$ and define

$$\boldsymbol{\theta}^{(k+1)} = \arg\max_{\boldsymbol{\theta}} Q(\boldsymbol{\theta}, \boldsymbol{\theta}^{(k)}) \tag{12}$$

It will turn out that we can maximize this quantity without numerical optimization. All we will need to do is solve a least-squares problem.

4. Iterate Step 2 and 3 until convergence.

Details. With a fixed parameter vector $\boldsymbol{\theta}^{(k)}$, the SDE (1) is specified completely, i.e., the drift and diffusion terms have no further unknowns.

Suppose we form R such time series. The expected log likelihood can then be approximated by

$$\begin{aligned} Q(\boldsymbol{\theta}, \boldsymbol{\theta}^{(k)}) &= \mathbb{E}_{\mathbf{z}|\mathbf{x},\boldsymbol{\theta}^{(k)}}[\log p(\mathbf{x}, \mathbf{z} \mid \boldsymbol{\theta})] \\ &\approx \frac{1}{R}\sum_{r=1}^{R}\Bigg[\sum_{j=1}^{N}\left[\sum_{i=1}^{d} -\frac{1}{2}\log(2\pi h \gamma_i^2)\right] \\ &\qquad - \frac{1}{2h}(y_j^{(r)} - y_{j-1}^{(r)} - h\sum_{k=1}^{M}\beta_k\phi_k(y_{j-1}^{(r)}))^T\Gamma^{-2}(y_j^{(r)} - y_{j-1}^{(r)} - h\sum_{\ell=1}^{M}\beta_\ell\phi_\ell(y_{j-1}^{(r)}))\Bigg] \end{aligned}$$

To maximize Q over $\boldsymbol{\theta}$, we first assume $\Gamma = \operatorname{diag}\gamma$ is known and maximize over β. This is a least squares problem. The solution is given by forming the matrix

$$\mathcal{M}_{k,\ell} = \frac{1}{R}\sum_{r=1}^{R}\sum_{j=1}^{N} h\phi_k^T(y_{j-1}^{(r)})\Gamma^{-2}\phi_\ell^T(y_{j-1}^{(r)})$$

and the vector

$$\rho_k = \frac{1}{R}\sum_{r=1}^{R}\sum_{j=1}^{N}\phi_k^T(y_{j-1}^{(r)})\Gamma^{-2}(y_j^{(r)} - y_{j-1}^{(r)}).$$

We then solve the system

$$\mathcal{M}\beta = \rho$$

for β. Now that we have β, we maximize Q over γ. The solution can be obtained in closed form:

$$\gamma_i^2 = \frac{1}{RNh}\sum_{r=1}^{R}\sum_{j=1}^{N}((y_j^{(r)} - y_{j-1}^{(r)} - h\sum_{\ell=1}^{M}\beta_\ell\phi_\ell(y_{j-1}^{(r)}))\cdot e_i)^2$$

where e_i is the i^{th} canonical basis vector in $\mathbb{R}^d$.

We demonstrate the method for 1, 2 and 3 dimensional systems.

- For the 1-dimensional system, we use the ? oscillator:

$$\mathrm{d}X(t) = (\alpha X(t) + \beta X(t)^2 + \gamma)\,\mathrm{d}t + g\,\mathrm{d}W(t) \tag{13}$$

- For the 2-dimensional system, we use the undamped Duffing oscillator:

$$\begin{aligned} \mathrm{d}X_1(t) &= X_2(t)\mathrm{d}t + g_1\,\mathrm{d}W_1(t) \\ \mathrm{d}X_2(t) &= (-X_1(t) - X_1^3(t))\mathrm{d}t + g_2\,\mathrm{d}W_2(t) \end{aligned}$$

- For the 3-dimensional case, we consider 2 different form of equations. The first one is the damped Duffing oscillator, a general form of the damped oscillator considered in the 2-dimensional case:

$$\begin{aligned} \mathrm{d}X_1(t) &= X_2(t)\,\mathrm{d}t + g_1\,\mathrm{d}W_1(t) \\ \mathrm{d}X_2(t) &= (\alpha X_1(t) - \beta X_1(t) - \delta X_2(t) + \gamma\cos(X_3(t)))\,\mathrm{d}t + g_2\,\mathrm{d}W_2(t) \\ \mathrm{d}X_3(t) &= \omega\,\mathrm{d}t + g_3\,\mathrm{d}W_3(t) \end{aligned}$$

- Another example considered for the 3-dimensional case is the Lorenz oscillator:

$$\begin{aligned}
\mathrm{d}X_1(t) &= \sigma(X_2(t) - X_1(t))\,\mathrm{d}t + g_1\,\mathrm{d}W_1(t)\\
\mathrm{d}X_2(t) &= (X_1(t)(\rho - X_3(t)))\mathrm{d}t + g_2\,\mathrm{d}W_2(t)\\
\mathrm{d}X_3(t) &= (X_1(t)X_2(t) - \beta X_3(t))\,\mathrm{d}t + g_3\,\mathrm{d}W_3(t)
\end{aligned}$$

For simplicity, consider the example where the $X \in \mathbb{R}^2$ and the highest degree of the Hermite polynomial is three, including four Hermite polynomials:

$$\begin{aligned}
f(x_1, x_2) &= \sum_{m=0}^{2} \sum_{i+j=0}^{i+j=m} \zeta_{i,j}\,\psi_{i,j}\\
&= \sum_{d=0}^{3} \sum_{i+j=0}^{i+j=3} \zeta_{i,j} H_i(x_1) H_j(x_2)\\
&= \sum_{i+j=0} \zeta_{i,j} H_i(x_1) H_j(x_2) + \sum_{i+j=1} \zeta_{i,j} H_i(x_1) H_j(x_2) + \sum_{i+j=2} \zeta_{i,j} H_i(x_1) H_j(x_2) + \sum_{i+j=3} \zeta_{i,j} H_i(x_1) H_j(x_2)\\
&= \zeta_{0,0} H_0(x_1) H_0(x_2) + \zeta_{0,1} H_0(x_1) H_1(x_2) + \zeta_{1,0} H_1(x_1) H_0(x_2) + \zeta_{0,2} H_0(x_1) H_2(x_2)\\
&+ \zeta_{2,0} H_2(x_1) H_0(x_2) + \zeta_{1,1} H_1(x_1) H_1(x_2) + \zeta_{0,3} H_0(x_1) H_3(x_2) + \zeta_{3,0} H_3(x_1) H_0(x_2)\\
&+ \zeta_{2,1} H_2(x_1) H_1(x_2) + \zeta_{1,2} H_1(x_1) H_2(x_2)
\end{aligned}$$

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