Learning Stochastic Dynamical Systems via Bridge Sampling

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Abstract

- The abstract paragraph should be indented ½ inch (3 picas) on both the left- and right-hand margins. Use 10 point type, with a vertical spacing (leading) of 11 points.

 The word **Abstract** must be centered, bold, and in point size 12. Two line spaces precede the abstract. The abstract must be limited to one paragraph.
- The goal of this work is to enable automatic discovery of stochastic differential equations (SDE) from
 time series data.
- o time series data.
 - Literature review.
 - 2. What is new and interesting about this work.
- 9 Points to cover:

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- data specification
 - Hermite polynomial and drift function representation
- Expectation and maximization formulas assuming data is filled in
- Filling data in with Brownian bridge
- MCMC iterations of brownian bridge using girsanov likelihood
 - 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, ...
- Note: constant noise case, not inferring the gvec

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 Idt. Let X_t denote an
- \mathbb{R}^d -valued stochastic process that evolves according to the Itô SDE

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

- ²⁴ For rigorous definitions of Brownian motion and SDE, see [Bhattacharya and Waymire 2009, Ok-
- sendal 2007]. 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

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

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

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

We see that the maximum degree of $H_{\alpha}(x)$ is $|\alpha|$. Hence we think of the double sum in (4) as first

$$f_j(x) = \sum_{i=1}^{\widetilde{M}} \beta_{(i)}^j H_{(i)}(x).$$
 (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 $\Delta t = \frac{1}{L} \sum_{j=1}^L (t_j - t_{j-1})$ be the average 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 . Note that $\widetilde{X}_{n+1} | \widetilde{X}_n = v$ has a $\mathcal{N}(v+f(v)h,h\Gamma^2)$ distribution. As h decreases, this Gaussian 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 [Roberts and Stramer, 2001; Papaspiliopoulos, Roberts, and Stramer, 2013]. Let us describe this method here. We suppose our current estimate of $\theta = (\beta, \gamma)$ is given. Then the goal is to generate a sample path of (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} . 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

$$d\hat{X}_t = \Gamma dW_t \tag{7}$$

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

$$\widehat{X}_{t} = \Gamma W_{t-t_{i}} + x_{i} - \frac{t - t_{i}}{t_{i+1} - t_{i}} (\Gamma W_{t_{i+1} - t_{i}} + x_{i} - x_{i+1})$$
(8)

Let $\mathbf{z}^{(r)}$ denote the r^{th} diffusion bridge sample path:

$$z^{(r)} \sim z \mid x, \beta^{(k)} \tag{9}$$

The observed and sampled data can be interleaved together to create a time series (completed data)

$$\mathbf{y}^{(r)} = \{y_j^{(r)}\}_{j=1}^N$$

of length N = LF + 1.

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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 $\theta = (\beta, \gamma)$.

- 1. Start with an initial guess for the parameters, $\theta^{(0)}$.
- 2. For the expectation (or E) step

$$Q(\boldsymbol{\theta}, \boldsymbol{\theta}^{(k)}) = \mathbb{E}_{\mathbf{z}|\mathbf{x}, \boldsymbol{\theta}^{(k)}}[\log p(\mathbf{x}, \mathbf{z} \mid \boldsymbol{\theta})]$$
(10)

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 θ and define

$$\boldsymbol{\theta}^{(k+1)} = \arg\max_{\boldsymbol{\theta}} Q(\boldsymbol{\theta}, \boldsymbol{\theta}^{(k)})$$
 (11)

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 $\theta^{(k)}$, 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

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

To maximize Q over θ , 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 .

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- We demonstrate the method for 1, 2 and 3 dimensional systems.
 - For the 1-dimensional system, we use the ? oscillator:

$$dX(t) = (\alpha X(t) + \beta X(t)^2 + \gamma) dt + g dW(t)$$
(12)

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

$$dX_1(t) = X_2(t)dt + g_1 dW_1(t)$$

$$dX_2(t) = (-X_1(t) - X_1^3(t))dt + g_2 dW_2(t)$$

• 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:

$$dX_1(t) = X_2(t) dt + g_1 dW_1(t)$$

$$dX_2(t) = (\alpha X_1(t) - \beta X_1(t) - \delta X_2(t) + \gamma \cos(X_3(t))) dt + g_2 dW_2(t)$$

$$dX_3(t) = \omega dt + g_3 dW_3(t)$$

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

$$dX_1(t) = \sigma(X_2(t) - X_1(t)) dt + g_1 dW_1(t)$$

$$dX_2(t) = (X_1(t)(\rho - X_3(t))) dt + g_2 dW_2(t)$$

$$dX_3(t) = (X_1(t)X_2(t) - \beta X_3(t)) dt + g_3 dW_3(t)$$

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:

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

92 **Expectation Maximization Steps**

- The data provided is in the form of a time series, $X \in \mathbb{R}^d$ at regular time points $t_l, 0 \le l \le L$. Note:
- 94 EM step in the sampling writeup.

95 3 Brownian bridge sampling

- 96 When the inter-observation time of the observed data X is large, the expectation step becomes less
- 97 accurate. To mitigate this problem, we can fill in the observed data with a Brownian bridge. We
- $_{98}$ generate many samples of the N-dimensional Brownian bridge and accept-reject samples using the
- 99 Metropolis-Hastings algorithm. The approximation of the likelihood is obtained using the Girsanov
- 100 likelihood function.

101 3.1 Brownian bridge

The \mathbb{R}^N dimensional Brownian bridge is defined by the integral:

$$I(t) = \int_0^t \frac{1-t}{1-T} \, dW(t)$$
 (13)

103 3.2 Metropolis Algorithm