Xinda Wu

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Content

- Generative Models and Variational Inference
 - ullet designs of variational distribution q_ϕ
- ② Deep Generative Models
 - deep latent Gaussian model
 - neural SDEs (inference, exact sampling, expressivity)
- Scalable Gradients for Neural SDEs
 - algorithm and experiment

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 (describes how data sets are generated, and by sampling from this model we can generate new data)

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- Bayesian Inference: $p_{\theta}(x|y) = p_{\theta}(y|x)p_{\theta}(x)/p_{\theta}(y)$
- Approximate Inference
 - MCMC
 - Variational Inference



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Also obtain the following inequality

$$-\log p_{\theta}(y) \leq \mathbf{F}_{\phi,\theta}(y)$$

'=' holds when $q_{\phi}(x|y) = p_{\theta}(x|y)$. The problem of MLE becomes:

$$\min_{\phi} \mathbf{F}_{\phi,\theta}(y)$$

$$\mathbf{F}_{\phi,\theta}(y) = \underbrace{D(q_{\phi}(x|y)||p_{\theta}(x))}_{\text{regulariser}} \underbrace{-\mathbf{E}_{q_{\phi}}[\log p_{\theta}(y|x_k)]}_{\text{reconstruction error}}$$

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Algorithm 2: Learning with Variational Inference
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Input: parameters \phi for variational distributions, \theta for the generative model (both initialized randomly). 

1 while the free energy F_{\phi,\theta} not converged do 

2 | y \leftarrow \{\text{Get mini-batch}\}

3 | compute the variational distribution q_{\phi}

4 | sample x \sim q_{\phi}(\cdot)

5 | compute the free energy F_{\phi,\theta}(y) \approx F_{\phi,\theta}(x,y)

6 | \Delta\theta \propto -\nabla_{\theta}F_{\phi,\theta}

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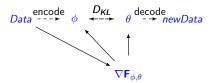
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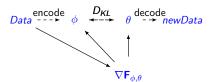
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Algorithm 4: Learning with Variational Inference

Input: parameters ϕ for variational distributions, θ for the generative model (both initialized randomly). 1 while the free energy $F_{\phi,\theta}$ not converged do $y \leftarrow \{\text{Get mini-batch}\}\$ compute the variational distribution q_{ϕ} sample $x \sim q_{\phi}(\cdot)$ compute the free energy $F_{\phi,\theta}(y) \approx F_{\phi,\theta}(x,y)$ $\Delta \theta \propto -\nabla_{\theta} F_{\phi,\theta}$ $\Delta \phi \propto -\nabla_{\phi} F_{\phi,\theta}$



Two problems in implementation:

- compute the gradients of free energy
- construct q_{ϕ} , balancing between richness and scalability

2

3

6

end

Mean-field Approximation (Parisi, 1988)

• assume the latent variables to be mutually independent and the distribution q_ϕ factorizes as follows

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- cannot capture the correlation between latent variables (fig.1)
- fail to fit a multimodal posterior (fig.2)
- Improvements: parameterize the correlation; use mixture model

$$q_{mix} = \sum_{i} \alpha_{i} q_{mf}^{(i)}$$



Figure 1: uncorrelated



Figure 2: not fitting multi-modal

Normalizing Flows (Rezende 2015)

 Idea: use a sequence of 'simple', differentiable, invertible transformations to construct an arbitrarily complex distribution

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• For each mapping $f: \mathbb{R}^d \to \mathbb{R}^d$, random variable x with distribution q(x), the resulting variable x' = f(x) has a distribution (chage of variable)

$$q(x') = q(x) \left| \det \frac{\partial f^{-1}}{\partial x'} \right| = q(x) \left| \det \frac{\partial f}{\partial x} \right|^{-1}.$$

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after K transformations

$$\ln q_K(x_K) = \ln q_0(x_0) - \sum_{i=1}^K \ln \left| \det \frac{\partial f_i}{\partial x_{i-1}} \right|$$

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• Law of the unconscious statistician (LOTUS): expectation w.r.t. q_K can be computed without explicitly knowing q_K

$$\mathbf{E}_{q_{K}}[g(x_{K})] = \mathbf{E}_{q_{0}}[g(f_{K} \circ ... \circ f_{1}(x_{0}))]$$



Planar flow:

$$f(x) = x + \mathbf{u}h(\mathbf{w}^{\mathsf{T}}x + b)$$

with $\mathbf{w} \in \mathbb{R}^d$, $\mathbf{u} \in \mathbb{R}^d$, $b \in \mathbb{R}$ parameters and $h : \mathbb{R} \to \mathbb{R}$ a smooth, element-wise nonlinear function

• Radial flow: $f(x) = x + \frac{\beta}{\alpha + r}(x - x_0)$ with $r = |x - x_0|$, $\alpha \in \mathbb{R}^+$, $\beta \in \mathbb{R}$

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- Free energy with planar flow

with $\psi(x) = h'(\mathbf{w}^{\mathsf{T}}x + b)\mathbf{w}$

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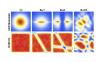


Figure 3: planarFlow



Figure 4: radialFlow



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Figure 5: approx posterior using planarFlow

Deep Latent Gaussian Model (Rezende, Kingma, 2014)

In this model, the latent variables $X_0,...,X_k$ and the observed variable Y are generated recursively according to

$$X_0 = Z_0$$

 $X_i = X_{i-1} + b_i(X_{i-1}) + \sigma_i Z_i, i = 1, ..., k$
 $Y \sim p(\cdot|X_k)$

where $Z_i \overset{i.i.d}{\sim} \mathcal{N}(0, I_d)$ in \mathbb{R}^d and $b_i : \mathbb{R}^d \to \mathbb{R}^d$ parametric nonlinear transformations, $\sigma_i \in \mathbb{R}^d \times \mathbb{R}^d$ sequence of matrices, $p(\cdot|\cdot)$ observation likelihood.

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- ullet representative power increases as $k \to \infty$. To avoid $n_{\theta} \to \infty$, consider its diffusion limit

and joint density

• Consider the continuous-time limit of DLGM, the latent object becomes a d-dimensional diffusion process

$$dX_t = b(X_t,t)dt + \sigma(X_t,t)dW_t, \ t \in [0,1]$$
 and the observed variable $Y \sim p(\cdot|X_1)$, latent space $\mathbb{W} = C([0,1];\mathbb{R}^d)$ and joint density

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and
$$b: \mathbb{R}^d imes [0,1] o \mathbb{R}^d$$

- ullet implemented by neural nets, heta weight parameters
- sufficiently well behaves (bounded, Lipschitz), admits a unique strong solution and a transition kernel



A Stochastic Control Problem

Consider the following stochastic control problem:

ullet controlled diffusion process $X^u = \{X^u_t\}_{t \in [0,1]}$ defined by

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Girsanov representation

X and X^u differ by a change of drift, by Girsanov formula

$$\frac{d\mathbf{P}^{u}}{d\mathbf{P}^{0}} = \exp\left(\int_{t}^{1} u_{s}^{T} dW_{s} + \frac{1}{2} \int_{t}^{1} \|u_{s}\|^{2} ds\right),$$

$$D(\mathbf{P}^{u} \| \mathbf{P}^{0}) = \mathbf{E}_{\mathbf{P}^{u}} \left[\log \frac{d\mathbf{P}^{u}}{d\mathbf{P}^{0}}\right] = \mathbf{E} \left[\frac{1}{2} \int_{t}^{1} \|u_{s}\|^{2} ds\right]$$

$$\Rightarrow J^{u}(x, t) := \mathbf{E} \left[\frac{1}{2} \int_{t}^{1} \|u_{s}\|^{2} ds - \log g(X_{1}^{u}) | X_{t}^{u} = x\right]$$

ullet Goal: to find the value function $v:\mathbb{R}^d imes [0,1] o \mathbb{R}_+$

$$v(x,t) := \inf_{u} J^{u}(x,t); \ v(\cdot,1) = -\log g(\cdot)$$

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• Bellman equation By the principle of optimality (Bellman, 1957), from t to t + dt,

$$v(x,t) = \min_{u} \left\{ v(x,t+dt) + \mathbf{E} \left[\frac{1}{2} \int_{t}^{t+dt} \|u_{s}\|^{2} ds \right] \right\}$$

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reduced to the following Cauchy problem

$$\frac{\partial v}{\partial t} + \mathcal{L}_t v = \frac{1}{2} \|\nabla v\|^2 \text{ on } \mathbb{R}^d \times [0,1]; v(\cdot,1) = -\log g(\cdot)$$

which can be solved by Feynman-Kac formula.



Theorem (Jamison, 1975; Dai Pra, 1991)

Consider the above control problem, then the value function is given by

$$v(x,t) = -\log \mathbf{E}[g(X_1)|X_t = x]$$

the optimal control is given by

$$u^*(x,t) = -\nabla v(x,t) = \nabla \log \mathbf{E}[g(X_1)|X_t = x]$$

the corresponding controlled diffusion X^{u^*} has the transition density

$$\kappa_{s,t}^*(x,y) = \kappa_{s,t}(x,y) \exp(v(x,s) - v(y,t))$$

where $\kappa_{s,t}(\cdot)$ is the transition density of uncontrolled process.

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• Variational upper-bound: (let g(x) = p(y|x) observation likelihood)

$$-\log \mathbf{E}[p(y|X_1)|X_0 = x] \le D(\mathbf{P^u}||\mathbf{P^0}) - \mathbf{E}[\log p(y|X_1^u)|X_0^u = x]$$

$$= \underbrace{\mathbf{E}\left[\frac{1}{2}\int_0^1 \|u_s\|^2 ds - \log p(y|X_1^u)\Big|X_0^u = x\right]}_{\mathbf{E^u}(y,y,\phi,\theta) := -}$$

Sampling problem: Given target distribution μ , and the prior process with $X_1 \sim \nu$

• $\exists u \text{ s.t. } X_1^u \sim \mu$?

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Theorem

Given a target μ at t=1, $X_0=0$, $X_1\sim \nu$, with $\mu\ll \nu$. Let $g=f=d\mu/d\nu$, then

$$X_1^* \sim \mu$$

and the optimal control u^* has the minimal energy

$$D(\mathbf{P^u}||\mathbf{P^0}) \ge D(d\mu||d\nu)$$

among all admissible controls u that induce the target distribution μ at t=1



Proof:

 $\bullet \ X_1^* \sim \mu$

$$\mathbb{P}[X_1^* \in A] = \int_A \kappa_{0,1}^*(0, y) dy$$

$$= \int_A \kappa_{0,1}(0, y) \exp(v(0, 0) - v(y, 1)) dy$$

$$= \int_A f d\nu$$

$$= \mu(A)$$

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$$= \mu(A)$$

Recall the entropy inequality

$$-\log\underbrace{\frac{\mathbf{E}[f(X_1)|X_0=0]}{\mathbf{E}[L_0]}}_{=\mathbf{E}_{\nu}[\frac{d\mu}{d\nu}]=1} \leq D(\mathbf{P^u}||\mathbf{P^0}) - \underbrace{\mathbf{E}[\log f(X_1^u)|X_0^u=0]}_{=\mathbf{E}_{\mu}[\log \frac{d\mu}{d\nu}]=D(d\mu||d\nu)}$$

$$\Longrightarrow D(\mathbf{P^u}||\mathbf{P^0}) \ge D(d\mu||d\nu)$$

$$\Longrightarrow \underbrace{D(d\mu||d\nu)}_{\text{minimal energy}} = \min_{u} \{ \frac{1}{2} \mathbf{E} [\int_{0}^{1} \|u_{s}\|^{2} ds] \}$$



Fix
$$b(x, t) \equiv 0$$
, $X_1 \sim \gamma_d$,

compute the value function

$$\begin{split} v(x,t) &= -\log \mathbf{E}[g(X_1)|X_t = x] = -\log \mathbf{E}[f(W_1)|X_t = x] \\ &= -\log \left((2\pi(1-t))^{-d/2} \int_t^1 f(y) \exp(-\frac{\|y - x\|^2}{2(1-t)}) dy \right) \\ &= -\log Q_{1-t} f(x) \end{split}$$

conpute the optimal control

$$u^*(x,t) = -\nabla v(x,t) = \underbrace{\nabla \log Q_{1-t}f(x)}_{ ext{F\"ollmer drift}}$$

Expressivity: could we approximate Föllmer drift by neural nets?

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- replace $Q_t f(x)$ by Monte Carlo estimate
- replace $f(\cdot)$ by a neural net approximation $\hat{f}(\cdot; \theta)$
- elementary operations(i.e. gradients,multiplications, reciprocals) can be computed by neural nets

$$\nabla \log Q_{1-t}f(x) \approx \nabla \log \left\{ \frac{1}{N} \sum_{n=1}^{N} \hat{f}(x + \sqrt{1-t}z_n; \theta) \right\}$$
$$= \frac{\sum_{n=1}^{N} \nabla \hat{f}(x + \sqrt{1-t}z_n; \theta)}{\sum_{n=1}^{N} \hat{f}(x + \sqrt{1-t}z_n; \theta)}$$

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- (regularity of the Föllmer drift) With assumption (A1), the Föllmer drift $b(x,t) = \nabla \log Q_{1-t}f(x)$ is bounded in the Euclidean norm,

$$\|\nabla \log Q_{1-t}f(x)\| \le \frac{L}{c}$$

for $x\in\mathbb{R}^d$, $t\in[0,1]$, and L the maximum of the Lipschitz constants of f and ∇f . And also, it is Lipschitz with Lipschitz constant $L/c+L^2/c^2$,

$$||b(x,t)-b(x',t)|| \le \left(\frac{L}{c} + \frac{L^2}{c^2}\right)||x-x'||.$$

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- (A2) activation function σ is differentiable and universal (any univariate Lipschitz function on a bounded interval can approximated well by a 2-layer MLP)
- (A3) both f and ∇f can be efficiently approximated by a neural net on any compact subset of \mathbb{R}^d



Theorem (Tzen 2019)

Suppose assumptions A1-A3 are in force. Let L denote the maximum of Lipschitz constants of f and ∇f . Then for any $0<\epsilon<16L^2/c^2$, there exists a neural net $\hat{v}:\mathbb{R}^d\times[0,1]\to\mathbb{R}^d$ with size polynomial in $1/\epsilon,d,L,c,1/c$, and the following holds:

If $\{\hat{X}_t\}_{t\in[0,1]}$ is the diffusion process governed by the Itô SDE

$$d\hat{X}_t = \hat{b}(\hat{X}_t, t)dt + dW_t, \, \hat{X}_0 = 0$$

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Proof:

Using probabilistic method and results from theory of empirical process to control the error incurred in each of the following steps:

- replace $Q_t f(x)$ by Monte Carlo estimate
- replace $f(\cdot)$ by a neural net approximation $\hat{f}(\cdot; \theta)$
- approximate the elementary operations by neural nets



Trainability

Consider the following problem of sensitivity analysis: (Gobet and Munos, 2005)

Given a d-dimensional Itô process,

$$dX_t^{\alpha} = b(X_t^{\alpha}, t; \alpha)dt + \sigma(X_t^{\alpha}, t; \alpha)dW_t; t \in [0, T]$$

where α is a n_{α} -dimensional parameter, and a function $\mathcal{L}: \mathbb{R}^d \to \mathbb{R}$ is given. We want to compute the gradient of the expectation

$$J(\alpha) := \mathbf{E}[\mathcal{L}(X_T^{\alpha})]$$

w.r.t α .

Note: for Neural SDE, $\alpha = (\phi, \theta)$, $\mathcal{L}(\cdot) = \mathbf{F}^u(y; \phi, \theta)(\cdot)$

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Path-wise forward

$$\nabla_{\alpha} J = \nabla_{\alpha} \mathbf{E}[\mathcal{L}(X_T^{\alpha})] = \mathbf{E}[\nabla \mathcal{L}(X_T^{\alpha}) \nabla_{\alpha} X_T^{\alpha}]$$

$$\frac{\partial X_t}{\partial \alpha_i} = \int_0^t \left(\frac{\partial b_s}{\partial \alpha_i} + \frac{\partial b_s}{\partial x} \frac{\partial X_s}{\partial \alpha_i} \right) ds + \sum_{l=1}^d \int_0^t \left(\frac{\partial \sigma_{s,l}}{\partial \alpha_i} + \frac{\partial \sigma_{s,l}}{\partial x} \frac{\partial X_s}{\partial \alpha_i} \right) dW_s^l$$

state sensitivities $\nabla_{\alpha} X_t$ will be computationally prohibitive as n_{α} increases largely



Idea:

• mimic the standard back-propagation of neural net

$$\frac{d\mathcal{L}}{dh_t} = \frac{d\mathcal{L}}{dh_{t+1}} \frac{dh_{t+1}}{dh_t}$$

- define the adjoint state $a_t = d\mathcal{L}/dx_t$
- derive a backward dynamic for the adjoint state which shows how the gradient propagate backward
- treating the parameters as additional part of the augmented state

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- treating the parameters as additional part of the augmented state

Advantages:

- adjoint state does not depend on parameters
- the gradients w.r.t parameters can be computed using adjoint state sensitivity
- ullet " n_lpha+1 systems " o " ${f 2}$ systems"

Some backward calculus: (Kunita, 2019)

• Two-sided filtration $\{\mathcal{F}_{s,t}\}_{s \leq t; s,t \in \mathbb{T}}$,

$$\mathcal{F}_{s,t} := \sigma(W_u - W_v : s \le u < v \le t), \, s,t \in \mathbb{T}(:=[0,T])$$

• Backward Wiener process $\{\check{W}_t\}_{t\in\mathbb{T}}$,

$$\check{W}_t = W_t - W_T, t \in \mathbb{T}$$

which is adapted to the backward filtration $\{\mathcal{F}_{t,T}\}_{t\in\mathbb{T}}$

• Backward Stratonovich integrals for continuous semimartingale \dot{Y}_t adapted to the backward filtration, define(in the L^2 sense)

$$\int_{s}^{T} \check{Y}_{t} \circ d\check{W}_{t} = \lim_{|\Pi| \to 0} \sum_{k=1}^{N} \frac{1}{2} \left(\check{Y}_{t_{k}} + \check{Y}_{t_{k-1}} \right) \left(\check{W}_{t_{k-1}} - \check{W}_{t_{k}} \right)$$

Stratonovich SDE

$$Z_T = z_0 + \int_0^T b(Z_t, t) dt + \sum_{i=1}^m \int_0^T \sigma_i(Z_t, t) \circ dW_t^{(i)}$$

with $b, \sigma \in C_b^{\infty,1}$ so that the SDE has unique strong solution.

• $\Phi_{s,t}(z) := Z_t^{s,z}$, the solution at time t when the process started at z at time s. Given a realization of the Wiener process, this defines

$$\mathcal{S} = \{\Phi_{s,t}\}_{s \leq t; s,t \in \mathbb{T}}$$

a collection of continuous maps from \mathbb{R}^d to itself, satisfying

$$\Phi_{s,u} = \Phi_{t,u} \circ \Phi_{s,t} \,, \ s < t < u$$

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- Stochastic flow of diffeomorphisms [Thm 3.7.1 (Kunita, 2019)]
 - ullet each $\Phi_{s,t}$ is a diffeomorphism $\mathbb{R}^d o \mathbb{R}^d$
 - $\check{\Psi}_{s,t} := \Phi_{s,t}^{-1}$, the *inverse flow* satisfies the backward SDE:

$$\check{\Psi}_{s,t}(z) = z - \int_s^t b(\check{\Psi}_{s,t}(z), u) du - \sum_{i=1}^m \int_s^t \sigma_i(\check{\Psi}_{s,t}(z), u) \circ d\check{W}_u^{(i)}$$

for $z \in \mathbb{R}^d$, $s, t \in \mathbb{T}$, $s \le t$.



Recall from
$$\frac{d\mathcal{L}}{dh_t} = \frac{d\mathcal{L}}{dh_{t+1}} \frac{dh_{t+1}}{dh_t}$$

•
$$A_{s,t}(z) := \nabla(\mathcal{L}(\Phi_{s,t}(z))) = \nabla\mathcal{L}(\Phi_{s,t}(z))\nabla\Phi_{s,t}(z)$$

$$\bullet \ \check{A}_{s,\mathit{T}}(z) := A_{s,\mathit{T}}(\check{\Psi}_{s,\mathit{T}}(z)) = \underbrace{\nabla \mathcal{L}(z)}_{\check{A}_{\mathit{T}}} \underbrace{\nabla \Phi_{s,\mathit{T}}(\check{\Psi}_{s,\mathit{T}}(z))}_{\partial Z_{\mathit{T}}/\partial Z_{s}}$$

• $J_{s,t}(z) := \nabla \check{\Psi}_{s,t}(z) = \partial Z_s/\partial Z_t$ satisfies,

$$J_{s,t}(z) = I_d - \int_s^t \nabla b(\check{\Psi}_{r,t}(z), r) J_{r,t}(z) dr$$
$$- \sum_{i=1}^m \int_s^t \nabla \sigma_i(\check{\Psi}_{r,t}(z), r) J_{r,t}(z) \circ d\check{W}_r^{(i)}$$

 $\bullet \ \, \textit{K}_{s,t}(\textit{z}) := \textit{J}_{s,t}(\textit{z})^{-1} = \nabla \Phi_{s,t}(\check{\Psi}_{s,t}(\textit{z})) = \partial \textit{Z}_t/\partial \textit{Z}_s \ \, \text{satisfies,}$

$$\begin{split} \mathcal{K}_{s,t}(z) &= I_d + \int_s^t \mathcal{K}_{r,t}(z) \nabla b(\check{\Psi}_{r,t}(z),r) dr \\ &+ \sum_{i=1}^m \int_s^t \mathcal{K}_{r,t}(z) \nabla \sigma_i(\check{\Psi}_{r,t}(z),r) \circ d\check{W}_r^{(i)} \end{split}$$



 $\bullet \ \, \check{A}_{s,T}(z) = \check{A}_T K_{s,T}(z) = \nabla \mathcal{L}(z) \nabla \Phi_{s,T}(\check{\Psi}_{s,T}(z))$

$$\check{A}_{s,T}(z) = \nabla \mathcal{L}(z) + \int_{s}^{T} \check{A}_{r,T}(z)^{\mathsf{T}} \nabla b(\check{\Psi}_{r,T}(z), r) dr
+ \sum_{i=1}^{m} \int_{s}^{T} \check{A}_{r,T}(z)^{\mathsf{T}} \nabla \sigma_{i}(\check{\Psi}_{r,T}(z), r) \circ d\check{W}_{r}^{(i)}$$

- consider augmented state $Y_t := (Z_t, \alpha)$, which satisfies a Stratonovich SDE with the drift function $\tilde{b}(y,t) = (b(z,t), \mathbf{0}_{n_{\alpha}})$ and the diffusion function $\tilde{\sigma}_i(y,t) = (\sigma_i(y,t), \mathbf{0}_{n_{\alpha}})$
- write the backward SDE for augmented adjoint $\check{A}^y := (\check{A}^z, \check{A}^\alpha)$ separately, we get the total gradient w.r.t parameters



Algorithm 5: Stochastic Ajoint Sensitivity (Stratonovich)

```
Input: parameters \alpha, start time t_0, stop time t_1, final state z_{t_1}, observation gradient \partial \mathcal{L}/z_{t_1}, drift function b(z,t,\alpha), diffusion function \sigma(z,t,\alpha), Wiener process sample path w(t).

1 def augmented drift \bar{b}(z_t,a_t,t,\alpha):
2 | return [-b(z_t,-t,\alpha), a_t^{\mathsf{T}}\partial b/\partial z, a_t^{\mathsf{T}}\partial b/\partial \alpha]
3 def augmented diffusion \bar{\sigma}(z_t,a_t,t,\alpha):
4 | return [-\sigma_i(z_t,-t,\alpha), a_t^{\mathsf{T}}\partial \sigma_i/\partial z, a_t^{\mathsf{T}}\partial \sigma_i/\partial \alpha]
5 def replicated noise \bar{w}(t):
6 | return [-w(-t),-w(-t),-w(-t)]
7 \begin{bmatrix} z_{t_0} \\ \partial \mathcal{L}/\partial z_{t_0} \\ \partial \mathcal{L}/\partial \alpha \end{bmatrix} = \mathrm{SDESolver} \begin{pmatrix} z_{t_1} \\ \partial \mathcal{L}/\partial z_{t_1} \\ 0_{n_{\alpha}} \end{bmatrix}, \bar{b}, \bar{\sigma}, \bar{w}, -t_1, -t_0
8 return \partial \mathcal{L}/\partial z_{t_0}, \partial \mathcal{L}/\partial \alpha
```

Algorithm 6: Stochastic Ajoint Sensitivity (Stratonovich)

Input: parameters α , start time t_0 , stop time t_1 , final state z_{t_1} , observation gradient $\partial \mathcal{L}/z_{t_1}$, drift function $b(z, t, \alpha)$, diffusion function $\sigma(z, t, \alpha)$, Wiener process sample path w(t).

```
1 def augmented drift \bar{b}(z_t, a_t, t, \alpha):
          return [-b(z_t, -t, \alpha), a_t^{\mathsf{T}} \partial b/\partial z, a_t^{\mathsf{T}} \partial b/\partial \alpha]
3 def augmented diffusion \bar{\sigma}(z_t, a_t, t, \alpha):
          return [-\sigma_i(z_t, -t, \alpha), a_t^{\mathsf{T}} \partial \sigma_i / \partial z, a_t^{\mathsf{T}} \partial \sigma_i / \partial \alpha]
5 def replicated noise \bar{w}(t):
          return [-w(-t), -w(-t), -w(-t)]
```

7
$$\begin{bmatrix} z_{t_0} \\ \partial \mathcal{L}/\partial z_{t_0} \\ \partial \mathcal{L}/\partial \alpha \end{bmatrix} = \text{SDESolver} \begin{pmatrix} z_{t_1} \\ \partial \mathcal{L}/\partial z_{t_1} \\ \mathbf{0}_{n_{\alpha}} \end{bmatrix}, \bar{b}, \bar{\sigma}, \bar{w}, -t_1, -t_0$$
8 return $\partial \mathcal{L}/\partial z_{t_1} = \partial \mathcal{L}/\partial \alpha$

8 return $\partial \mathcal{L}/\partial z_{t_0}$, $\partial \mathcal{L}/\partial \alpha$

Table 1: L denotes the numbers of steps in the SDE solving. n_{α} is the dimension of the parameter, and d is the dimension of the system state.

Method	Memory	Time
Path-wise Forward Sensitivity Adjoint Sensitivity	$\mathcal{O}(1)$ $\mathcal{O}(1)$	$\frac{\mathcal{O}(L\cdot(n_{\alpha}+d))}{\mathcal{O}(L)}$

Implementation

• parameterize the prior and the approximate posterior using SDEs

$$\begin{split} dZ_t &= b_{\theta}(Z_t,t)dt + \sigma(Z_t,t)dW_t \quad \text{(prior)} \\ d\tilde{Z}_t &= \tilde{b}_{\phi}(\tilde{Z}_t,t)dt + \sigma(\tilde{Z}_t,t)dW_t \quad \text{(approx.post.)} \end{split}$$

sharing the same diffusion function σ , with $Z_0 = \tilde{Z}_0 = z_0 \in \mathbb{R}^d$, and $u : \mathbb{R}^d \times [0, T] \to \mathbb{R}^m$ satisfies $b_{\sigma}(z, t) - b_{\theta}(z, t) = \sigma(z, t)u(z, t)$.

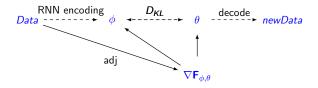
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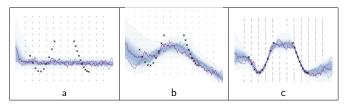
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Modelling time series: use a recurrent neural net (RNN) to encode the posterior process,
 MLP to implement the prior process and KL divergence penalty for parameter update

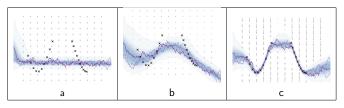


• Presetting a stochastic process, and using simulation data to play neural SDEs

- Presetting a stochastic process, and using simulation data to play neural SDEs
- fitting to a Ornstein-Uhlenbeck process

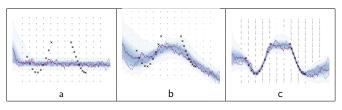


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- Possible issues in practice:
 - SDE solvers may not be working well
 - overfitting
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- Possible issues in practice:
 - SDE solvers may not be working well
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 - KL divergence may not be well-defined or constantly 0, when we deal with high dimensional data supported in low dimensional sub-manifold
- Unofficial package from Google Research: https://github.com/google-research/torchsde



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Thank you!

