

Lecture 15 - Sampling and Diffusion Models

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- ① 问题背景
- ② 从给定概率密度中采样
 - MCMC
 - Langevin算法
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 - Score-Matching Langevin Dynamics (SMLD)
 - Denoising Diffusion Probabilistic Models
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 - More Topics

采样问题

采样(Sampling)

设 μ 是一个概率分布，**采样问题**是指：如何获得随机样本 X ，使得 X 的分布为 μ 。

- 计算：Monte Carlo 方法；
- 优化：模拟退火；
- 生成式AI；
- ...

生成

生成任务

在机器学习中，**生成任务**是指模型的目标为生成新的数据实例，这些实例与已有数据(训练数据)具有相似的特征或模式，常见的生成任务包括文本生成、图像生成、视频生成等。

- 从给定概率密度中采样：

给定函数 $H : \mathcal{X}^d \rightarrow \mathbb{R}$ ，如何获得服从概率分布 $\mu(x) \propto e^{-H(x)}$ 的随机样本？

- 从给定数据中采样：

给定数据样本 $X_1, X_2, \dots, X_n \sim p_{data}$ ，如何生成服从数据分布 p_{data} 的数据样本？

目录

1 问题背景

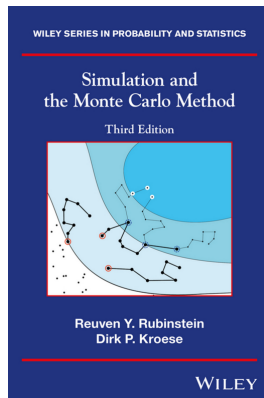
2 从给定概率密度中采样

- MCMC
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随机模拟



- 随机数的模拟
- 基本随机变量的模拟
- 随机过程/随机向量的模拟

Markov Chain Monte Carlo (MCMC)

定义 1 (Markov链)

一个 **Markov链** 是一个随机过程 $\{X_n\}_{n=0}^{\infty}$, 满足: 未来状态只依赖于当前状态, 而与过去状态无关。即, 对于任意的 n 和状态 i, j , 有:

$$P(X_{n+1} = j \mid X_n = i, X_{n-1} = x_{n-1}, \dots, X_0 = x_0) = P(X_{n+1} = j \mid X_n = i)$$

进一步, 如果转移概率不随时间变化, 即:

$$P(X_{n+1} = j \mid X_n = i) = P(X_1 = j \mid X_0 = i) \quad \forall n$$

那么我们称 **Markov链** 是**时齐的**。

Markov Chain Monte Carlo (MCMC)

定义 2 (平稳分布)

设 $\{X_n\}_{n=0}^{\infty}$ 是一个 Markov 链，其状态空间为 S 。如果存在一个概率分布 π 满足以下条件：

$$\pi(j) = \sum_{i \in S} \pi(i)P(i, j) \quad \forall j \in S$$

其中 $P(i, j)$ 是从状态 i 转移到状态 j 的概率，则称 π 为该 Markov 链的平稳分布 (Stationary Distribution) / 不变测度 (Invariant Measure)。

MCMC 的思想即是：构造一个 Markov 链，使得它的平稳分布是我们采样的目标分布 π 。那么从任意状态分布 (容易获得样本的分布) 出发，经过充分的状态转移，就能获得目标分布 π 的样本。

Metropolis-Hastings 算法

- 1 选择初始状态 x_0 。
- 2 对于每一步 $n = 1, 2, \dots, N$:
 - 从提议分布 $q(x'|x_{n-1})$ 中生成候选状态 x' 。
 - 计算接受概率:

$$\alpha = \min \left(1, \frac{\pi(x')q(x_{n-1}|x')}{\pi(x_{n-1})q(x'|x_{n-1})} \right)$$

- 以概率 α 接受候选状态 x' , 否则保持状态 x_{n-1} 。

设 X_{n-1} 是当前状态, X_n 是下一个状态, Metropolis-Hastings 算法转移概率可以表示为:

$$P(X_n = x' | X_{n-1} = x) = q(x'|x) \cdot \min \left(1, \frac{\pi(x')q(x|x')}{\pi(x)q(x'|x)} \right)$$

容易check Metropolis-Hastings 算法的平稳分布是目标分布 π

$$\pi(x) = \sum_y \pi(y)q(x|y) \cdot \min \left(1, \frac{\pi(x)q(y|x)}{\pi(y)q(x|y)} \right)$$

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Langevin Dynamics

定义 3 (Langevin Dynamics)

给定势能函数 (potential) $V(x)$, *Langevin Dynamics* 是如下形式的 SDE

$$dX_t = -\nabla V(X_t)dt + \sqrt{2}dB_t. \quad (1)$$

其解一般称之为 *Langevin Diffusion*.

Langevin 扩散的生成元为

$$\mathcal{L}_{LD}f = -\nabla V \cdot \nabla f + \Delta f$$

生成元的伴随为

$$\mathcal{L}_{LD}^*g = \nabla \cdot (g\nabla V) + \Delta g$$

Langevin Dynamics

Kolmogorov backward方程:

$$\frac{\partial}{\partial t} P_t f(x) = \mathcal{L}_{LD} P_t f(x) = -\nabla V(x) \cdot \nabla P_t f(x) + \Delta P_t f(x)$$

Fokker-Planck方程:

$$\partial_t \mu(x, t) = \mathcal{L}_{LD}^* \mu(x, t) = \nabla \cdot (\mu(x, t) \nabla V(x)) + \Delta \mu(x, t)$$

命题 1

Langevin 扩散 $dX_t = -\nabla V(X_t)dt + \sqrt{2}dB_t$ 的不变测度为

$$d\pi(x) \propto e^{-V(x)} dx$$

Langevin Algorithm

定义 4 (Langevin Algorithm)

对Langevin扩散Euler-Maruyama离散化:

$$X_{(k+1)h} := X_{kh} - h\nabla V(X_{kh}) + \sqrt{2}(B_{(k+1)h} - B_{kh}).$$

我们得到了一种Langevin扩散的实现方式, 称为(Unadjusted) Langevin Algorithm, ULA, 或者 Langevin Monte Carlo, LMC. 其中 h 是迭代步长, k 是迭代轮数。

由于时间离散化, Langevin Monte Carlo与Langevin Dynamics不再一致, Langevin Monte Carlo的平稳分布也不再是目标分布。一般可以通过Metropolis调整保证采样分布的准确性。

Metropolis-adjusted Langevin Algorithm (MALA)

MALA

- Proposal step: same as in ULA

$$Z_{k+1} = X_k - h\nabla V(X_k) + \sqrt{2h}\xi_k$$

- Accept-reject step: go to

$$X_{k+1} = \begin{cases} Z_{k+1} & \text{with probability } \min \left\{ 1, \frac{\pi(Z_{k+1})\mathcal{P}_{Z_{k+1}}(X_k)}{\pi(X_k)\mathcal{P}_{X_k}(Z_{k+1})} \right\} \\ X_k & \text{with the remaining probability.} \end{cases}$$

注意到给定 X_k , 提议分布是一个均值为 $X_k - h\nabla V(X_k)$, 方差为 $2h\mathbb{I}_n$ 的高斯分布, 即提议分布显式表达为

$$\mathcal{P}_z(x) = \frac{1}{(2\pi \cdot 2h)^{\frac{n}{2}}} \exp \left(-\frac{\|x - (z - h\nabla V(z))\|_2^2}{4h} \right).$$

接受概率也具有显式表达:

$$\min \left\{ 1, \exp \left(-V(z) - \frac{1}{4h} \|x - (z - h\nabla V(z))\|_2^2 + V(x) + \frac{1}{4h} \|z - (x - h\nabla V(x))\|_2^2 \right) \right\}$$

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耦合方法

定义 5 (Wasserstein distance)

概率测度 μ 和 ν 之间的2-Wasserstein Distance定义为:

$$W_2(\mu, \nu) := \inf_{\gamma \in \mathcal{C}(\mu, \nu)} \left(\int \|x - y\|^2 \gamma(\mathrm{d}x, \mathrm{d}y) \right)^{\frac{1}{2}}. \quad (2)$$

其中 $\mathcal{C}(\mu, \nu)$ 是 μ 和 ν 的耦合(Couplings)构成的空间, $\|\cdot\|$ 是欧式范数。

定理 1

设 $\{X_t\}$ 为初值为 $X_0 \sim \mu_0$, 平稳分布为 $\mu \propto e^{-V}$ 的Langevin扩散, 假设 μ 是 α -强log-concave的, 那么

$$W_2^2(\mu_t, \mu) \leq \exp(-2\alpha t) W_2^2(\mu_0, \mu).$$

耦合方法

定理 2

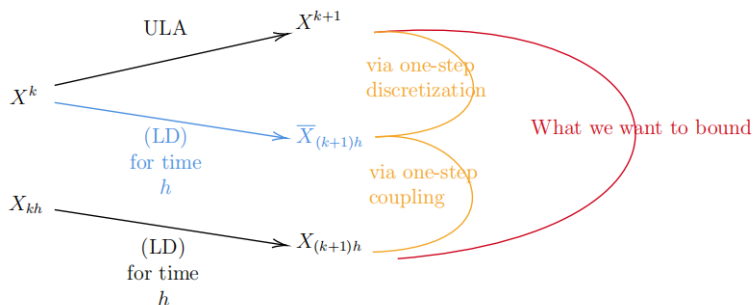
对于 $k \in \mathbb{N}$, 记 μ_{kh} 为 LMC 的第 k 轮迭代的分布, $h > 0$ 为迭代步长。设目标分布为 $\mu \propto e^{-V}$, 满足 $\alpha I_d \leq \nabla^2 V \leq \beta I_d$. 如果 $h \lesssim \frac{1}{\beta\kappa}$, 那么对于所有 $N \in \mathbb{N}$,

$$W_2(\mu_{Nh}, \mu) \leq \exp\left(-\frac{\alpha Nh}{2}\right) W_2(\mu_0, \mu) + O\left(\frac{\beta d^{1/2} h^{1/2}}{\alpha}\right).$$

如果 $h = O(\frac{\varepsilon^2}{\beta\kappa d})$, 那么对于任意 $\varepsilon \in [0, \sqrt{d}]$, 在

$$N = O\left(\frac{\kappa^2 d}{\varepsilon^2} \log \frac{\sqrt{\alpha} W_2(\mu_0, \mu)}{\varepsilon}\right)$$

轮迭代之后, 有 $\sqrt{\alpha} W_2(\mu_{Nh}, \mu) \leq \varepsilon$.

Proof sketch¹

- 计算ULA和LD之间的一步时间离散化误差；
- 通过LD在 W_2 距离下的指数压缩性分析多步迭代误差。

¹<https://chewisinho.github.io/main.pdf> Sec.4.1

泛函不等式方法²

耦合方法的分析要求目标分布的log-concavity，泛函不等式方法可以减弱这一假设。

定义 6 (Log-Sobolev inequality)

称 ν 满足 α Log-Sobolev inequality，如果对于 $\mathbb{E}_\nu[g^2] < \infty$ 的光滑函数 $g : \mathbb{R}^n \rightarrow \mathbb{R}$,

$$\text{Ent}_\nu(g) \triangleq \mathbb{E}_\nu[g^2 \log g^2] - \mathbb{E}_\nu[g^2] \log \mathbb{E}_\nu[g^2] \leq \frac{2}{\alpha} \mathbb{E}_\nu[\|\nabla g\|^2]. \quad (3)$$

定义 7 (Poincaré inequality)

称 ν 满足 α Poincaré inequality，如果对于光滑函数 $g : \mathbb{R}^n \rightarrow \mathbb{R}$, 有

$$\text{Var}_\nu(g) \triangleq \mathbb{E}_\nu[g^2] - \mathbb{E}_\nu[g]^2 \leq \frac{1}{\alpha} \mathbb{E}_\nu[\|\nabla g\|^2]. \quad (4)$$

²Bakry D, Gentil I, Ledoux M. Analysis and geometry of Markov diffusion operators[M]. Cham: Springer, 2014.

KL散度

定义 8 (KL散度)

μ 对于 ν 的KL散度定义为

$$\text{KL}(\mu\|\nu) = H_\nu(\mu) = \int_{\mathbb{R}^n} \mu(x) \log \frac{\mu(x)}{\nu(x)} dx. \quad (5)$$

命题 2 (Pinsker's inequality)

$$d_{\text{TV}}(\mu, \nu)^2 \leq \frac{1}{2} H_\nu(\mu).$$

命题 3 (Talagrand inequality)

若 ν 满足 α Log-Sobolev inequality

$$\frac{\alpha}{2} W_2(\mu, \nu)^2 \leq H_\nu(\mu).$$

Pinsker's inequality和Talagrand inequality说明KL散度是一个相对更强的距离度量，我们bound KL散度自然能够给出TV和 W_2 距离的界。

KL散度+LSI下Langevin Dynamics的指数收敛性

定理 3

若 $\nu \propto e^{-V}$ 满足 α LSI, 那么Langevin Dynamics

$$dX_t = -\nabla V(X_t)dt + \sqrt{2}dW_t$$

的分布 μ_t 满足:

$$H_\nu(\mu_t) \leq e^{-2\alpha t} H_\nu(\mu_0).$$

进一步地, $W_2(\mu_t, \nu) \leq \sqrt{\frac{2}{\alpha} H_\nu(\mu_0)} e^{-\alpha t}.$

Proof sketch³

定义 9 (Fisher information)

μ 对于 ν 的Fisher information定义为

$$J_\nu(\mu) = \int_{\mathbb{R}^n} \mu(x) \left\| \nabla \log \frac{\mu(x)}{\nu(x)} \right\|^2 dx. \quad (6)$$

令 $g^2 = \frac{\mu}{\nu}$, Log-Sobolev inequality可以得到KL散度和Fisher information的如下关系:

$$H_\nu(\mu) \leq \frac{1}{2\alpha} J_\nu(\mu).$$

³Vempala S, Wibisono A. Rapid convergence of the unadjusted langevin algorithm: Isoperimetry suffices[J]. Advances in neural information processing systems, 2019, 32.

Proof sketch

引理 1

分布 μ_t 满足:

$$\frac{d}{dt} H_\nu(\mu_t) = -J_\nu(\mu_t). \quad (7)$$

利用Langevin Dynamics的Fokker-Planck方程: $\partial_t \mu_t = \nabla \cdot (\mu_t \nabla V(x)) + \Delta \mu_t$ 计算可得。

由Log-Sobolev inequality,

$$H_\nu(\mu) \leq \frac{1}{2\alpha} J_\nu(\mu).$$

结合(7)式, 有

$$\frac{d}{dt} H_\nu(\mu_t) \leq -2\alpha H_\nu(\mu_t)$$

两边积分有

$$H_\nu(\mu_t) \leq e^{-2\alpha t} H_\nu(\mu_0).$$

KL散度+LSI下LMC的指数收敛性⁴

定理 4

若 $\nu := e^{-V}$ 满足 α LSI 并且是 L -smooth的($-LI \preceq \nabla^2 V(x) \preceq LI$ for all $x \in \mathbb{R}^n$), 那么对于任意 $x_0 \sim \mu_0$ 满足 $H_\nu(\mu_0) < \infty$, 步长 $0 < \eta \leq \frac{\alpha}{4L^2}$ 的ULA

$$x_{k+1} = x_k - \eta \nabla V(x_k) + \sqrt{2\eta} z_k$$

的分布 $x_k \sim \mu_k$ 满足:

$$H_\nu(\mu_k) \leq e^{-\alpha\eta k} H_\nu(\mu_0) + \frac{8\eta n L^2}{\alpha}.$$

因此, 对任意精度 $\delta > 0$, 为了 $H_\nu(\mu_k) < \delta$, LMC需要满足步长 $\eta \leq \frac{\alpha}{4L^2} \min\{1, \frac{\delta}{4n}\}$, 并且经过 $k \geq \frac{1}{\alpha\eta} \log \frac{2H_\nu(\mu_0)}{\delta}$ 次迭代。

⁴Vempala S, Wibisono A. Rapid convergence of the unadjusted langevin algorithm: Isoperimetry suffices[J]. Advances in neural information processing systems, 2019, 32.

Proof sketch

- 给出一步LMC迭代的界

引理 2

若 $\nu := e^{-V}$ 满足 α Log-Sobolev inequality 并且 L -smooth, 步长 $0 < \eta \leq \frac{\alpha}{4L^2}$, 那么LMC满足:

$$H_\nu(\mu_{k+1}) \leq e^{-\alpha\eta} H_\nu(\mu_k) + 6\eta^2 n L^2.$$

$$\frac{d}{dt} H_\nu(\mu_t) \leq -\frac{3}{4} J_\nu(\mu_t) + \frac{4t^2 L^4}{\alpha} H_\nu(\mu_0) + 2t^2 n L^3 + 2tn L^2.$$

由Log-Sobolev inequality

$$\frac{d}{dt} H_\nu(\mu_t) \leq -\frac{3\alpha}{2} H_\nu(\mu_t) + \frac{4t^2 L^4}{\alpha} H_\nu(\mu_0) + 2t^2 n L^3 + 2tn L^2.$$

$t = 0$ 到 $t = \eta$ 积分, 整理可得引理2。

- 给出多步迭代的界

Rényi散度+PI下Langevin Dynamics的指数收敛性

定义 10 (Rényi散度)

对于 $q > 0, q \neq 1$, 概率分布 μ 对于 ν 的 q 阶Rényi散度定义为:

$$R_{q,\nu}(\mu) := \frac{1}{q-1} \log F_{q,\nu}(\mu), \quad (8)$$

其中

$$F_{q,\nu}(\mu) := \mathbb{E}_{\nu} \left[\left(\frac{\mu}{\nu} \right)^q \right] = \int_{\mathbb{R}^n} \nu(x) \frac{\mu(x)^q}{\nu(x)^q} dx = \int_{\mathbb{R}^n} \frac{\mu(x)^q}{\nu(x)^{q-1}} dx.$$

Rényi散度来源于Rényi熵: $H_q(\mu) := \frac{1}{q-1} \log \int \mu(x)^q dx$.

定理 5

若 $\nu := e^{-f}$ 满足 α Poincaré inequality, $q \geq 2$, 那么Langevin Dynamics的分布 μ_t 满足:

$$R_{q,\nu}(\mu_t) \leq \begin{cases} R_{q,\nu}(\mu_0) - \frac{2\alpha t}{q} & \text{if } R_{q,\nu}(\mu_0) \geq 1 \text{ and as long as } R_{q,\nu}(\mu_t) \geq 1, \\ e^{-\frac{2\alpha t}{q}} R_{q,\nu}(\mu_0) & \text{if } R_{q,\nu}(\mu_0) \leq 1. \end{cases}$$

Rényi散度+PI下LMC的指数收敛性

定理 6

若 ν_η 满足 β Poincaré inequality, $q \geq 1$, $\nu := e^{-V}$ 是 L -smooth的,
且 $1 \leq R_{2q, \nu_\eta}(\mu_0) < \infty$, 令 $0 < \eta \leq \min \left\{ \frac{1}{3L}, \frac{1}{9\beta} \right\}$, $q > 1$, 那么对于 $k \geq k_0 := \frac{2q}{\beta\eta} (R_{2q, \nu_\eta}(\mu_0) - 1)$, LMC满足:

$$R_{q, \nu}(\mu_k) \leq \left(\frac{q - \frac{1}{2}}{q - 1} \right) e^{-\frac{\beta\eta(k-k_0)}{2q}} + R_{2q-1, \nu}(\nu_\eta).$$

对任意精度 $\delta > 0$, 为了 $R_{q, \nu}(\mu_k) \leq \delta$, LMC需要满足步长 $\eta = \Theta \left(\min \left\{ \frac{1}{L}, \gamma_{2q-1} \left(\frac{\delta}{2} \right) \right\} \right)$, 其中 $\gamma_q(\delta) = \sup \{ \eta > 0: R_{q, \nu}(\nu_\eta) \leq \delta \}$, 并且经过 $k = \Theta \left(\frac{1}{\beta\eta} (R_{2q, \nu_\eta}(\mu_0) + \log \frac{1}{\delta}) \right)$ 次迭代。

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生成——从数据分布中采样

数据分布是**未知的**，我们仅有一些样本，扩散模型的基本思想是：

(1) 学习数据分布；(2) 根据学习到的数据分布生成实例

SMLD

Recall that Langevin dynamics

$$dX_t = -\nabla V(X_t)dt + \sqrt{2}dB_t$$

具有不变测度 $\pi \propto e^{-V}$ ，若我们需要采样 p_{data} ，可以通过

$$dX_t = \nabla \log p_{data}(X_t)dt + \sqrt{2}dB_t.$$

其中 $\nabla \log p$ 称为概率分布 p 的 Score function.

如果 p_{data} 已知，那么可以显式计算 Score，然而在生成任务中，我们需要从数据中学习 Score $\nabla \log p_{data}$ ，通过神经网络近似

$$\min_{\theta} \mathbb{E}[\|\nabla \log p_{data}(X) - s_{\theta}(X)\|_2^2]$$

其中 s_{θ} 为参数为 θ 的神经网络。

Score Matching

• Score matching⁵

$$\begin{aligned} & \mathbb{E}_{X \sim p_{data}} \|\nabla \log p_{data}(X) - s_{\theta}(X)\|_2^2 \\ &= \underbrace{\mathbb{E} \|\nabla \log p_{data}(X)\|_2^2}_{\text{does not depend on } \theta} - 2\mathbb{E} \langle s_{\theta}(X), \nabla \log p_{data}(X) \rangle + \mathbb{E} \|s_{\theta}(X)\|_2^2. \end{aligned}$$

计算第二项

$$\begin{aligned} -\mathbb{E} \langle s_{\theta}(X), \nabla \log p_{data}(X) \rangle &= -\int \langle s_{\theta}(x), \nabla \log p_{data}(x) \rangle p_{data}(x) dx \\ &= \int \nabla \cdot s_{\theta}(x) p_{data}(x) dx = \mathbb{E} \nabla \cdot s_{\theta}(X), \end{aligned}$$

那么我们的优化问题即为:

$$\min_{\theta} \mathbb{E}_{X \sim p_{data}} [\|s_{\theta}(X)\|_2^2 + 2\nabla \cdot s_{\theta}(X)].$$

⁵Hyvärinen A, Dayan P. Estimation of non-normalized statistical models by score matching[J]. Journal of Machine Learning Research, 2005, 6(4).

Denoising score matching

实际训练神经网络优化经验风险函数：

$$\min_{\theta} \frac{1}{N} \sum_{i=1}^N [\|s_{\theta}(x_i)\|_2^2 + 2\nabla \cdot s_{\theta}(x_i)] .$$

然而高维情形计算散度项 $\nabla \cdot s_{\theta}(x_i)$ 比较困难，考虑通过 **Denoising score matching** 避免散度的计算。

- Denoising score matching⁶

考虑扰动 $\tilde{x} = x + \sigma z$ ，其中 $z \sim N(0, I)$ ，Denoising score matching 的目标为

$$\min_{\theta} \mathbb{E}_{q_{\sigma}(\tilde{x}|x)p_{data}(x)} [\|s_{\theta}(\tilde{x}) - \nabla_{\tilde{x}} \log q_{\sigma}(\tilde{x} | x)\|_2^2] .$$

可以证明 $s_{\theta^*}(\tilde{x}) = \nabla_{\tilde{x}} \log q_{\sigma}(\tilde{x})$ 几乎处处成立⁷，其中 $q_{\sigma}(\tilde{x}) \triangleq \int q_{\sigma}(\tilde{x} | x)p_{data}(x)dx$ 。

⁶Yang Song and Stefano Ermon. Generative modeling by estimating gradients of the data distribution. In Advances in Neural Information Processing Systems, pp. 11895–11907, 2019.

⁷Vincent P. A connection between score matching and denoising autoencoders[J]. Neural computation, 2011, 23(7): 1661-1674.

Denoising score matching

虽然只有在 σ 比较小的时候，有

$$s_{\theta^*}(\tilde{x}) = \nabla_{\tilde{x}} \log q_{\sigma}(\tilde{x}) \approx \nabla_x \log p_{\text{data}}(x)$$

但是 $q_{\sigma}(\tilde{x} | x) \sim \mathcal{N}(x, \sigma^2 I)$ 是条件高斯的在计算上十分高效

$$\begin{aligned} \nabla_{\tilde{x}} \log q_{\sigma}(\tilde{x} | x) &= \nabla_{\tilde{x}} \log \frac{1}{(\sqrt{2\pi\sigma^2})^d} \exp \left\{ -\frac{\|\tilde{x} - x\|^2}{2\sigma^2} \right\} \\ &= \nabla_{\tilde{x}} \left\{ -\frac{\|\tilde{x} - x\|^2}{2\sigma^2} - \log(\sqrt{2\pi\sigma^2})^d \right\} \\ &= -\frac{\tilde{x} - x}{\sigma^2}. \end{aligned}$$

Denoising score matching的优化问题即为：

$$\min_{\theta} \mathbb{E}_{x \sim p_{\text{data}}, \tilde{x} \sim \mathcal{N}(x, \sigma^2 I)} \left\| s_{\theta}(\tilde{x}, \sigma) + \frac{\tilde{x} - x}{\sigma^2} \right\|_2^2. \quad (9)$$

Noise Conditional Score Networks

(9)中参数化的神经网络模型 $s_{\theta}(x, \sigma)$ 称为 Noise Conditional Score Networks。由于只有在 σ 比较小的时候, 有

$$s_{\theta^*}(x, \sigma) \approx \nabla_x \log p_{\text{data}}(x)$$

考虑设计一个time schedule, 使得 $\sigma_t \rightarrow 0$.

Algorithm 1 Annealed Langevin dynamics.

Require: $\{\sigma_i\}_{i=1}^L, \epsilon, T$.

```

1: Initialize  $\tilde{\mathbf{x}}_0$ 
2: for  $i \leftarrow 1$  to  $L$  do
3:    $\alpha_i \leftarrow \epsilon \cdot \sigma_i^2 / \sigma_L^2$   $\triangleright \alpha_i$  is the step size.
4:   for  $t \leftarrow 1$  to  $T$  do
5:     Draw  $\mathbf{z}_t \sim \mathcal{N}(0, I)$ 
6:      $\tilde{\mathbf{x}}_t \leftarrow \tilde{\mathbf{x}}_{t-1} + \frac{\alpha_i}{2} \mathbf{s}_{\theta}(\tilde{\mathbf{x}}_{t-1}, \sigma_i) + \sqrt{\alpha_i} \mathbf{z}_t$ 
7:   end for
8:    $\tilde{\mathbf{x}}_0 \leftarrow \tilde{\mathbf{x}}_T$ 
9: end for
return  $\tilde{\mathbf{x}}_T$ 

```

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- MCMC
- Langevin算法
- Langevin算法的收敛性分析

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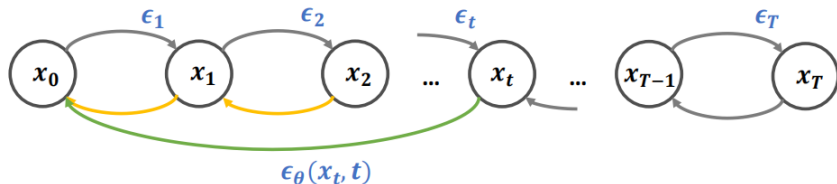
- Score-Matching Langevin Dynamics (SMLD)
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Denoising Diffusion Probabilistic Models, DDPM⁸

Forward/Diffusion Process



Reverse/Denoise Process



⁸Denoising diffusion probabilistic models. Advances in Neural Information Processing Systems, 33, 2020

Forward Diffusion Process

- 一步加噪过程

$$q(\mathbf{x}_t | \mathbf{x}_{t-1}) = \mathcal{N}(\mathbf{x}_t; \sqrt{\alpha_t} \mathbf{x}_{t-1}, (1 - \alpha_t) \mathbf{I})$$

$$\mathbf{x}_t = \sqrt{\alpha_t} \mathbf{x}_{t-1} + \sqrt{(1 - \alpha_t)} \boldsymbol{\epsilon}_{t-1}, \quad \text{where } \boldsymbol{\epsilon}_{t-1} \sim \mathcal{N}(0, \mathbf{I}).$$

- t 步加噪过程

命题 4

条件分布 $q(\mathbf{x}_t | \mathbf{x}_0)$ 为

$$q(\mathbf{x}_t | \mathbf{x}_0) = \mathcal{N}(\mathbf{x}_t; \sqrt{\bar{\alpha}_t} \mathbf{x}_0, (1 - \bar{\alpha}_t) \mathbf{I}),$$

其中 $\bar{\alpha}_t = \prod_{i=1}^t \alpha_i$. 即 $\mathbf{x}_t = \sqrt{\bar{\alpha}_t} \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \boldsymbol{\epsilon}_0$.

能够计算 $q(\mathbf{x}_t | \mathbf{x}_0)$ 的好处在于给定 \mathbf{x}_0 ，给一个 t 可以直接得到 \mathbf{x}_t .

Proof

$$\begin{aligned}
 \mathbf{x}_t &= \sqrt{\alpha_t} \mathbf{x}_{t-1} + \sqrt{1 - \alpha_t} \epsilon_{t-1} \\
 &= \sqrt{\alpha_t} (\sqrt{\alpha_{t-1}} \mathbf{x}_{t-2} + \sqrt{1 - \alpha_{t-1}} \epsilon_{t-2}) + \sqrt{1 - \alpha_t} \epsilon_{t-1} \\
 &= \sqrt{\alpha_t \alpha_{t-1}} \mathbf{x}_{t-2} + \underbrace{\sqrt{\alpha_t} \sqrt{1 - \alpha_{t-1}} \epsilon_{t-2} + \sqrt{1 - \alpha_t} \epsilon_{t-1}}_{\mathbf{w}_1}.
 \end{aligned}$$

由于 ϵ_{t-2} 和 ϵ_{t-1} 都是标准高斯的， \mathbf{w}_1 是均值为0的高斯，我们下面计算协方差

$$\begin{aligned}
 \mathbb{E}[\mathbf{w}_1 \mathbf{w}_1^T] &= [(\sqrt{\alpha_t} \sqrt{1 - \alpha_{t-1}})^2 + (\sqrt{1 - \alpha_t})^2] \mathbf{I} \\
 &= [\alpha_t (1 - \alpha_{t-1}) + 1 - \alpha_t] \mathbf{I} = [1 - \alpha_t \alpha_{t-1}] \mathbf{I}.
 \end{aligned}$$

延用记号 ϵ_t

$$\begin{aligned}
 \mathbf{x}_t &= \sqrt{\alpha_t \alpha_{t-1}} \mathbf{x}_{t-2} + \sqrt{1 - \alpha_t \alpha_{t-1}} \epsilon_{t-2} \\
 &= \sqrt{\alpha_t \alpha_{t-1} \alpha_{t-2}} \mathbf{x}_{t-3} + \sqrt{1 - \alpha_t \alpha_{t-1} \alpha_{t-2}} \epsilon_{t-3} \\
 &= \cdots = \sqrt{\prod_{i=1}^t \alpha_i} \mathbf{x}_0 + \sqrt{1 - \prod_{i=1}^t \alpha_i} \epsilon_0.
 \end{aligned}$$

Reverse Denoising Process

我们希望用一个神经网络实现降噪过程，即

$$p_{\theta}(\mathbf{x}_{t-1}|\mathbf{x}_t) \approx q(\mathbf{x}_{t-1}|\mathbf{x}_t)$$

由Markov性，

$$q(\mathbf{x}_t|\mathbf{x}_{t-1}) = \frac{q(\mathbf{x}_{t-1}|\mathbf{x}_t)q(\mathbf{x}_t)}{q(\mathbf{x}_{t-1})} \xrightarrow{\text{condition on } \mathbf{x}_0} q(\mathbf{x}_t|\mathbf{x}_{t-1}, \mathbf{x}_0) = \frac{q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0)q(\mathbf{x}_t|\mathbf{x}_0)}{q(\mathbf{x}_{t-1}|\mathbf{x}_0)}$$

在优化神经网络的过程中转化为⁹¹⁰

$$p_{\theta}(\mathbf{x}_{t-1}|\mathbf{x}_t) \approx q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0)$$

⁹Luo C. Understanding diffusion models: A unified perspective[J]. arXiv preprint arXiv:2208.11970, 2022.

¹⁰Chan S H. Tutorial on Diffusion Models for Imaging and Vision[J]. arXiv preprint arXiv:2403.18103, 2024.

Reverse Denoising Process

命题 5

条件分布 $q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0)$ 为一个高斯分布 $\mathcal{N}(\mathbf{x}_{t-1}; \boldsymbol{\mu}_q(\mathbf{x}_t, \mathbf{x}_0), \boldsymbol{\Sigma}_q(t))$, 其中

$$\boldsymbol{\mu}_q(\mathbf{x}_t, \mathbf{x}_0) = \frac{(1 - \bar{\alpha}_{t-1})\sqrt{\alpha_t}}{1 - \bar{\alpha}_t} \mathbf{x}_t + \frac{(1 - \alpha_t)\sqrt{\bar{\alpha}_{t-1}}}{1 - \bar{\alpha}_t} \mathbf{x}_0$$

$$\boldsymbol{\Sigma}_q(t) = \frac{(1 - \alpha_t)(1 - \sqrt{\bar{\alpha}_{t-1}})}{1 - \bar{\alpha}_t} \mathbf{I}$$

$$\begin{aligned} q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0) &= \frac{q(\mathbf{x}_t|\mathbf{x}_{t-1}, \mathbf{x}_0)q(\mathbf{x}_{t-1}|\mathbf{x}_0)}{q(\mathbf{x}_t|\mathbf{x}_0)} \\ &= \frac{\mathcal{N}(\mathbf{x}_t; \sqrt{\alpha_t}\mathbf{x}_{t-1}, (1 - \alpha_t)\mathbf{I})\mathcal{N}(\mathbf{x}_{t-1}; \sqrt{\bar{\alpha}_{t-1}}\mathbf{x}_0, (1 - \bar{\alpha}_{t-1})\mathbf{I})}{\mathcal{N}(\mathbf{x}_t; \sqrt{\bar{\alpha}_t}\mathbf{x}_0, (1 - \bar{\alpha}_t)\mathbf{I})} \\ &\propto \exp \left\{ - \left[\frac{(\mathbf{x}_t - \sqrt{\alpha_t}\mathbf{x}_{t-1})^2}{2(1 - \alpha_t)} + \frac{(\mathbf{x}_{t-1} - \sqrt{\bar{\alpha}_{t-1}}\mathbf{x}_0)^2}{2(1 - \bar{\alpha}_{t-1})} - \frac{(\mathbf{x}_t - \sqrt{\bar{\alpha}_t}\mathbf{x}_0)^2}{2(1 - \bar{\alpha}_t)} \right] \right\} \end{aligned}$$

Reverse Denoising Process

注意到, 给定加噪schedule, $\Sigma_q(t)$ 是已知的, 所以我们只需要参数化均值部分, 即

$$p_{\theta}(x_{t-1}|x_t) = \mathcal{N}(x_{t-1}; \mu_{\theta}, \Sigma_q(t))$$

两个高斯分布之间的KL散度可以容易计算:

$$D_{\text{KL}}(q(x_{t-1}|x_t, x_0) \parallel p_{\theta}(x_{t-1}|x_t)) = \frac{1}{2\sigma_q^2(t)} \left[\|\mu_{\theta} - \mu_q\|_2^2 \right]$$

注意到

$$\begin{aligned} \mu_q(x_t, x_0) &= \frac{(1 - \bar{\alpha}_{t-1})\sqrt{\alpha_t}}{1 - \bar{\alpha}_t} x_t + \frac{(1 - \alpha_t)\sqrt{\bar{\alpha}_{t-1}}}{1 - \bar{\alpha}_t} x_0 \\ &= \frac{(1 - \bar{\alpha}_{t-1})\sqrt{\alpha_t}}{1 - \bar{\alpha}_t} x_t + \frac{(1 - \alpha_t)\sqrt{\bar{\alpha}_{t-1}}}{1 - \bar{\alpha}_t} \frac{x_t - \sqrt{1 - \bar{\alpha}_t} \epsilon_0}{\sqrt{\bar{\alpha}_t}} \\ &= \frac{1}{\sqrt{\alpha_t}} x_t - \frac{1 - \alpha_t}{\sqrt{1 - \bar{\alpha}_t} \sqrt{\alpha_t}} \epsilon_0 \end{aligned}$$

Denoising Diffusion Probabilistic Models

考虑

$$\mu_{\theta}(\mathbf{x}_t, t) = \frac{1}{\sqrt{\alpha_t}} \mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1 - \bar{\alpha}_t} \sqrt{\alpha_t}} \epsilon_{\theta}(\mathbf{x}_t, t)$$

我们要学习的目标其实是一个 Denoiser $\epsilon_{\theta}(\mathbf{x}_t, t)$ 。

Algorithm 1 Training

```

1: repeat
2:    $\mathbf{x}_0 \sim q(\mathbf{x}_0)$ 
3:    $t \sim \text{Uniform}(\{1, \dots, T\})$ 
4:    $\epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ 
5:   Take gradient descent step on
        $\nabla_{\theta} \|\epsilon - \epsilon_{\theta}(\sqrt{\bar{\alpha}_t} \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \epsilon, t)\|^2$ 
6: until converged
  
```

Algorithm 2 Sampling

```

1:  $\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ 
2: for  $t = T, \dots, 1$  do
3:    $\mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$  if  $t > 1$ , else  $\mathbf{z} = \mathbf{0}$ 
4:    $\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left( \mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1 - \bar{\alpha}_t}} \epsilon_{\theta}(\mathbf{x}_t, t) \right) + \sigma_t \mathbf{z}$ 
5: end for
6: return  $\mathbf{x}_0$ 
  
```

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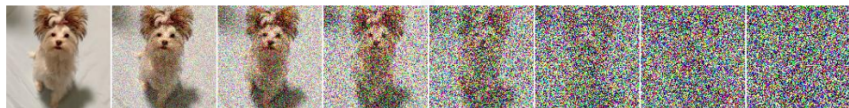
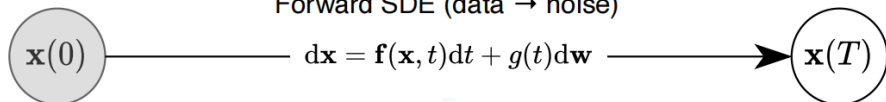
1 问题背景

2 从给定概率密度中采样

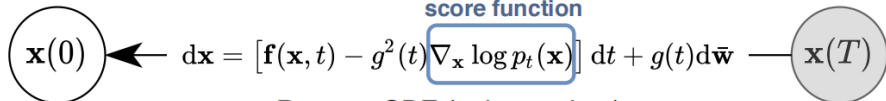
- MCMC
- Langevin算法
- Langevin算法的收敛性分析

3 从给定数据中采样：扩散模型

- Score-Matching Langevin Dynamics (SMLD)
- Denoising Diffusion Probabilistic Models
- **Score-based Generative Models**
- More Topics

Score-based Generative Models, SGM¹¹Forward SDE (data \rightarrow noise)

score function

Reverse SDE (noise \rightarrow data)

¹¹Y. Song, J. Sohl-Dickstein, D. P. Kingma, A. Kumar, S. Ermon, and B. Poole. Score-based generative modeling through stochastic differential equations. In Proc. ICLR, 2021.

Reverse SDE

定理 7

对于如下SDE:

$$d\mathbf{x} = \mathbf{f}(\mathbf{x}, t)dt + \mathbf{G}(\mathbf{x}, t)d\mathbf{w}, \quad (10)$$

它的 *Reverse SDE* 为

$$d\mathbf{x} = \{\mathbf{f}(\mathbf{x}, t) - \nabla \cdot [\mathbf{G}(\mathbf{x}, t)\mathbf{G}(\mathbf{x}, t)^\top] - \mathbf{G}(\mathbf{x}, t)\mathbf{G}(\mathbf{x}, t)^\top \nabla_{\mathbf{x}} \log p_t(\mathbf{x})\}dt + \mathbf{G}(\mathbf{x}, t)d\bar{\mathbf{w}}$$

Proof Sketch

SDE (10)的Fokker-Planck方程为

$$\begin{aligned}\frac{\partial p_t(\mathbf{x})}{\partial t} &= - \sum_{i=1}^d \frac{\partial}{\partial x_i} [f_i(\mathbf{x}, t) p_t(\mathbf{x})] + \frac{1}{2} \sum_{i=1}^d \sum_{j=1}^d \frac{\partial^2}{\partial x_i \partial x_j} \left[\sum_{k=1}^d G_{ik}(\mathbf{x}, t) G_{jk}(\mathbf{x}, t) p_t(\mathbf{x}) \right] \\ &= - \sum_{i=1}^d \frac{\partial}{\partial x_i} [f_i(\mathbf{x}, t) p_t(\mathbf{x})] + \frac{1}{2} \sum_{i=1}^d \frac{\partial}{\partial x_i} \left[\sum_{j=1}^d \frac{\partial}{\partial x_j} \left[\sum_{k=1}^d G_{ik}(\mathbf{x}, t) G_{jk}(\mathbf{x}, t) p_t(\mathbf{x}) \right] \right].\end{aligned}$$

注意到

$$\begin{aligned}& \sum_{j=1}^d \frac{\partial}{\partial x_j} \left[\sum_{k=1}^d G_{ik}(\mathbf{x}, t) G_{jk}(\mathbf{x}, t) p_t(\mathbf{x}) \right] \\ &= \sum_{j=1}^d \frac{\partial}{\partial x_j} \left[\sum_{k=1}^d G_{ik}(\mathbf{x}, t) G_{jk}(\mathbf{x}, t) \right] p_t(\mathbf{x}) + \sum_{j=1}^d \sum_{k=1}^d G_{ik}(\mathbf{x}, t) G_{jk}(\mathbf{x}, t) p_t(\mathbf{x}) \frac{\partial}{\partial x_j} \log p_t(\mathbf{x}) \\ &= p_t(\mathbf{x}) \nabla \cdot [\mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^\top] + p_t(\mathbf{x}) \mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^\top \nabla_{\mathbf{x}} \log p_t(\mathbf{x})\end{aligned}$$

Proof Sketch

回代Fokker-Planck方程

$$\begin{aligned}
 \frac{\partial p_t(\mathbf{x})}{\partial t} &= - \sum_{i=1}^d \frac{\partial}{\partial x_i} [f_i(\mathbf{x}, t) p_t(\mathbf{x})] \\
 &+ \frac{1}{2} \sum_{i=1}^d \frac{\partial}{\partial x_i} \left[p_t(\mathbf{x}) \nabla \cdot [\mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^\top] + p_t(\mathbf{x}) \mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^\top \nabla_{\mathbf{x}} \log p_t(\mathbf{x}) \right] \\
 &= - \sum_{i=1}^d \frac{\partial}{\partial x_i} \left\{ f_i(\mathbf{x}, t) p_t(\mathbf{x}) \right. \\
 &\quad \left. - \frac{1}{2} \left[\nabla \cdot [\mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^\top] + \mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^\top \nabla_{\mathbf{x}} \log p_t(\mathbf{x}) \right] p_t(\mathbf{x}) \right\} \\
 &\triangleq - \sum_{i=1}^d \frac{\partial}{\partial x_i} [\tilde{f}_i(\mathbf{x}, t) p_t(\mathbf{x})],
 \end{aligned}$$

做时间逆转

$$\frac{\partial p_t(\mathbf{x})}{\partial t} = - \sum_{i=1}^d \frac{\partial}{\partial x_i} [-\tilde{f}_i(\mathbf{x}, t) p_t(\mathbf{x})] \tag{11}$$

Proof Sketch

整理(11), 得

$$\frac{\partial p_t(\mathbf{x})}{\partial t} = - \sum_{i=1}^d \frac{\partial}{\partial x_i} [\bar{f}_i(\mathbf{x}, t) p_t(\mathbf{x})] + \frac{1}{2} \sum_{i=1}^d \sum_{j=1}^d \frac{\partial^2}{\partial x_i \partial x_j} \left[\sum_{k=1}^d G_{ik}(\mathbf{x}, t) G_{jk}(\mathbf{x}, t) p_t(\mathbf{x}) \right]$$

其中

$$\bar{\mathbf{f}}(\mathbf{x}, t) = \mathbf{f}(\mathbf{x}, t) - \nabla \cdot [\mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^{\top}] - \mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^{\top} \nabla_{\mathbf{x}} \log p_t(\mathbf{x})$$

所以Reverse SDE 为

$$d\mathbf{x} = \{\mathbf{f}(\mathbf{x}, t) - \nabla \cdot [\mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^{\top}] - \mathbf{G}(\mathbf{x}, t) \mathbf{G}(\mathbf{x}, t)^{\top} \nabla_{\mathbf{x}} \log p_t(\mathbf{x})\} dt + \mathbf{G}(\mathbf{x}, t) d\bar{\mathbf{w}}$$

Forward Process of DDPM & OU Process

考虑离散时间 $i = 1, 2, \dots, N$, DDPM的前向加噪过程

$$\mathbf{x}_i = \sqrt{1 - \beta_i} \mathbf{x}_{i-1} + \sqrt{\beta_i} \mathbf{z}_{i-1}, \quad \mathbf{z}_{i-1} \sim \mathcal{N}(0, \mathbf{I}).$$

定义时间步长 $\Delta t = \frac{1}{N}$, $t \in \{0, 1, \dots, \frac{N-1}{N}\}$ 。加噪schedule为

$$\beta_i = \beta \left(\frac{i}{N} \right) \cdot \frac{1}{N} = \beta(t + \Delta t) \Delta t, \quad N \rightarrow \infty, \beta \left(\frac{i}{N} \right) \rightarrow \beta(t)$$

于是

$$\begin{aligned} \mathbf{x}(t + \Delta t) &= \sqrt{1 - \beta(t + \Delta t) \Delta t} \mathbf{x}(t) + \sqrt{\beta(t + \Delta t) \Delta t} \mathbf{z}(t) \\ &\approx \mathbf{x}(t) - \frac{1}{2} \beta(t + \Delta t) \Delta t \mathbf{x}(t) + \sqrt{\beta(t + \Delta t) \Delta t} \mathbf{z}(t) \\ &\approx \mathbf{x}(t) - \frac{1}{2} \beta(t) \Delta t \mathbf{x}(t) + \sqrt{\beta(t) \Delta t} \mathbf{z}(t), \end{aligned}$$

当 $\Delta t \rightarrow 0$,

$$d\mathbf{x} = -\frac{1}{2} \beta(t) \mathbf{x} dt + \sqrt{\beta(t)} d\mathbf{w}.$$

Denoiser和Score的联系

引理 3 (Tweedie Formula)

对于一个高斯随机变量 $z \sim \mathcal{N}(z; \mu_z, \Sigma_z)$, 有

$$\mathbb{E}[\mu_z | z] = z + \Sigma_z \nabla_z \log p(z)$$

在DDPM中, 我们证明过

$$q(x_t | x_0) = \mathcal{N}(x_t; \sqrt{\bar{\alpha}_t} x_0, (1 - \bar{\alpha}_t) \mathbf{I})$$

应用Tweedie Formula

$$\mathbb{E}[\mu_{x_t} | x_t] = \sqrt{\bar{\alpha}_t} x_0 = x_t + (1 - \bar{\alpha}_t) \nabla \log p(x_t)$$

带入到 $\mu_q(x_t, x_0) = \frac{\sqrt{\alpha_t}(1 - \bar{\alpha}_{t-1})x_t + \sqrt{\bar{\alpha}_{t-1}}(1 - \alpha_t)x_0}{1 - \bar{\alpha}_t}$ 中计算, 有

$$\mu_q(x_t, x_0) = \frac{1}{\sqrt{\alpha_t}} x_t + \frac{1 - \alpha_t}{\sqrt{\alpha_t}} \nabla \log p(x_t)$$

Denoiser和Score的联系

可以通过学习到的Score计算

$$\mu_{\theta}(\mathbf{x}_t, t) = \frac{1}{\sqrt{\alpha_t}} \mathbf{x}_t + \frac{1 - \alpha_t}{\sqrt{\alpha_t}} \mathbf{s}_{\theta}(\mathbf{x}_t, t)$$

又由

$$\mathbf{x}_0 = \frac{\mathbf{x}_t + (1 - \bar{\alpha}_t) \nabla \log p(\mathbf{x}_t)}{\sqrt{\bar{\alpha}_t}} = \frac{\mathbf{x}_t - \sqrt{1 - \bar{\alpha}_t} \boldsymbol{\epsilon}_0}{\sqrt{\bar{\alpha}_t}}$$

可以得到Denoiser和Score的联系

$$\nabla \log p(\mathbf{x}_t) = -\frac{1}{\sqrt{1 - \bar{\alpha}_t}} \boldsymbol{\epsilon}_0$$

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Reverse OU Process

- Forward process

$$d\mathbf{X}_t = -\beta_t \mathbf{X}_t dt + \sqrt{2\beta_t} d\mathbf{B}_t$$

- Backward process (BP)

$$d\mathbf{Y}_t = \beta_{T-t} \{\mathbf{Y}_t + 2\nabla \log p_{T-t}(\mathbf{Y}_t)\} dt + \sqrt{2\beta_{T-t}} d\mathbf{B}_t$$

Diffusion Model 收敛性分析

- Girsanov 定理
 - Chen S, et al. Sampling is as easy as learning the score: theory for diffusion models with minimal data assumptions. ICLR. 2023.
 - Chen H, et al. Improved analysis of score-based generative modeling: User-friendly bounds under minimal smoothness assumptions. ICML. 2023.
 - Benton J, et al. Nearly d -Linear Convergence Bounds for Diffusion Models via Stochastic Localization. ICLR. 2024
- Log-Sobolev inequality
 - Convergence for score-based generative modeling with polynomial complexity. NeuralIPS. 2022.
 - Convergence of score-based generative modeling for general data distributions. 2023.
- 其他
 - A Note on the Convergence of Denoising Diffusion Probabilistic Models. TMLR. 2024.

Reverse Diffusion Monte Carlo¹²

设采样目标为 $x \propto e^{-f_*(x)}$, 考虑 Reverse Diffusion Process

$$d\mathbf{X}_t = \beta_{T-t} \{ \mathbf{X}_t + 2\nabla \log p_{T-t}(\mathbf{X}_t) \} dt + \sqrt{2\beta_{T-t}} d\mathbf{B}_t$$

引理 4

The score function can be rewritten as

$$\nabla_{\mathbf{x}} \log p_{T-t}(\mathbf{x}) = \mathbb{E}_{\mathbf{x}_0 \sim q_{T-t}(\cdot|\mathbf{x})} \frac{e^{-(T-t)} \mathbf{x}_0 - \mathbf{x}}{(1 - e^{-2(T-t)})},$$

where

$$q_{T-t}(\mathbf{x}_0|\mathbf{x}) \propto \exp \left(-f_*(\mathbf{x}_0) - \frac{\|\mathbf{x} - e^{-(T-t)} \mathbf{x}_0\|^2}{2(1 - e^{-2(T-t)})} \right).$$

¹²Huang X, Dong H, Yifan H A O, et al. Reverse diffusion monte carlo[C]//The Twelfth International Conference on Learning Representations. 2024.

Reverse Diffusion Monte Carlo

Algorithm 1 RDMC: reverse diffusion Monte Carlo

1: **Input:** Initial particle $\tilde{\mathbf{x}}_0$ sampled from \tilde{p}_0 , Terminal time T , Step size η, η' , Sample size n .

2: **for** $k = 0$ to $\lfloor T/\eta \rfloor - 1$ **do**

3: Set $\mathbf{v}_k = \mathbf{0}$;

4: Create n Monte Carlo samples to estimate

$$\mathbf{v}_k \approx \mathbb{E}_{\mathbf{x} \sim q_{T-t}} \left[-\frac{\tilde{\mathbf{x}}_{k\eta} - e^{-(T-k\eta)} \mathbf{x}}{(1 - e^{-2(T-k\eta)})} \right], \text{ where } q_{T-t}(\mathbf{x} | \tilde{\mathbf{x}}_{k\eta}) \propto \exp \left(-f_*(\mathbf{x}) - \frac{\|\tilde{\mathbf{x}}_{k\eta} - e^{-(T-k\eta)} \mathbf{x}\|^2}{2(1 - e^{-2(T-k\eta)})} \right).$$

5: $\tilde{\mathbf{x}}_{(k+1)\eta} = e^\eta \tilde{\mathbf{x}}_{k\eta} + (e^\eta - 1) \mathbf{v}_k + \xi$ where ξ is sampled from $\mathcal{N}(0, (e^{2\eta} - 1) \mathbf{I}_d)$.

6: **end for**

7: **Return:** $\tilde{\mathbf{x}}_{\lfloor T/\eta \rfloor \eta}$.
