

# Deep Learning lecture 7

## Normalizing Flow & Score-Based Model

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Mar-31

# Today's Topic

- Normalizing Flow
  - A generative model class that has the best sampling and likelihood properties
- Score-based generative model
  - A different framework to tackle general energy-based models

# Today's Topic

- **Normalizing Flow**
  - A generative model class that has the best sampling and likelihood properties
- Score-based generative model
  - A different framework to tackle general energy-based models

# Latent Variable Model (Recap)

- $p(x, z) = p(z)p(x|z)$ 
  - Given  $z$ , we use  $p(x|z)$  to generate  $x$  for a consistent probability distribution
  - Hard to estimate the exact likelihood  $p(x) = \int_z p(x|z)p(z)$
  - Variational inference by ELBO
- **Can we simplify the generation process?**

# A Simplified Generation Process

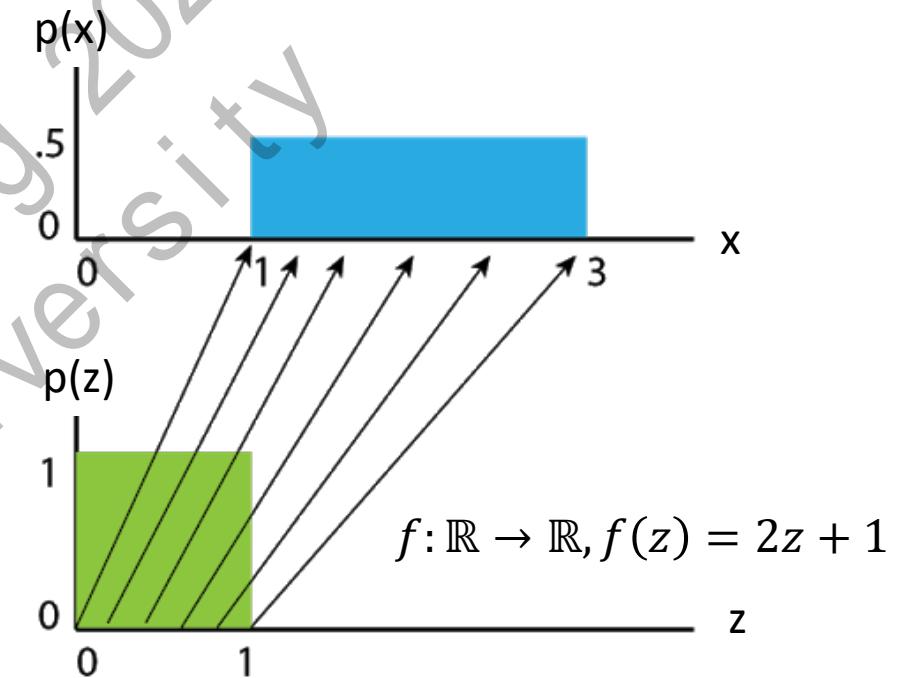
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- Can we simplify the generation process?
  - We can directly apply a deterministic process  $f: z \rightarrow x$
  - E.g.,  $z \sim N(0, I)$  and  $x = f(z)$ 
    - $x \sim N(\mu, \sigma^2)$  is equivalent to  $z \sim N(0, 1)$ ,  $x = \mu + \sigma \cdot z$

# A Simplified Generation Process

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  - E.g.,  $z \sim N(0, I)$  and  $x = f(z)$
- **Can we make the likelihood  $p(x)$  tractable?**
  - So that we can directly run MLE for training ...

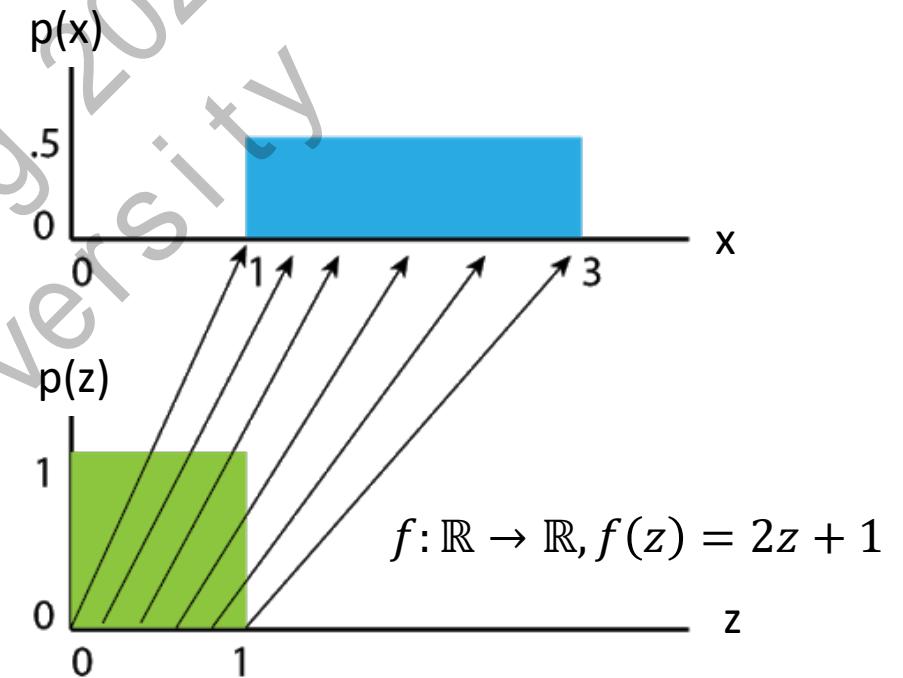
# 1-D Example

- Goal: design  $x = f(z; \theta)$  s.t.
  - Assume  $z$  is from an “easy” distribution
  - $p(x) = p(f(z; \theta))$  has a tractable likelihood
- Uniform:  $z \sim \text{unif}(0,1)$ 
  - Density  $p(z) = 1$
  - $x = 2z + 1$ , then  $p(x) = ?$



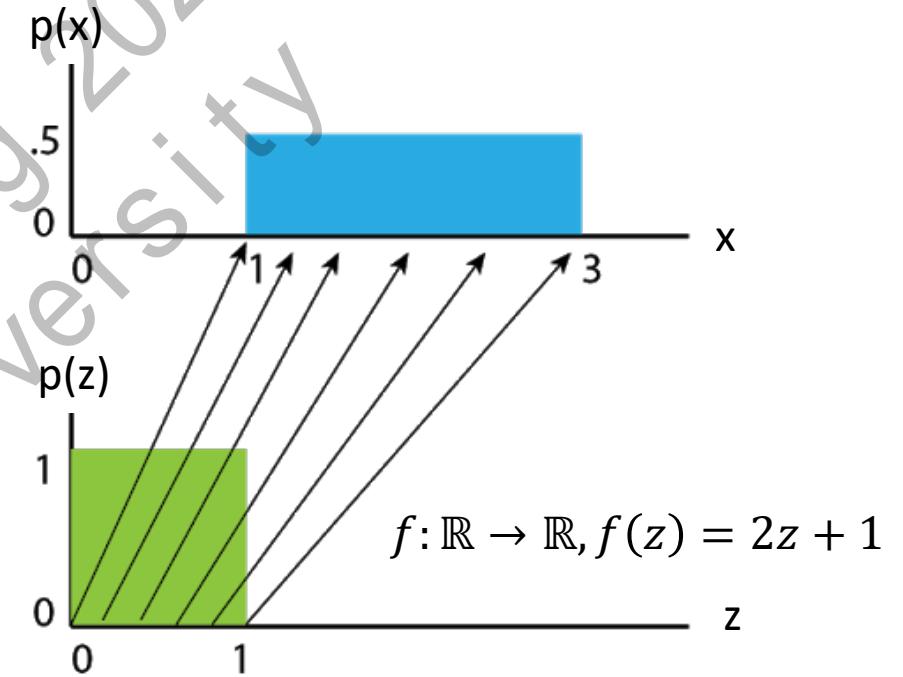
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  - $x = 2z + 1$ , then  $p(x) = \frac{1}{2}$ 
    - $x = a \cdot z + b$ , then  $p(x) = 1/|a|$  (for  $a \neq 0$ )
  - General 1-D case:  $x = f(z)$ ,  $p(x) = ?$ 
    - Assume  $f(z)$  is a bijection



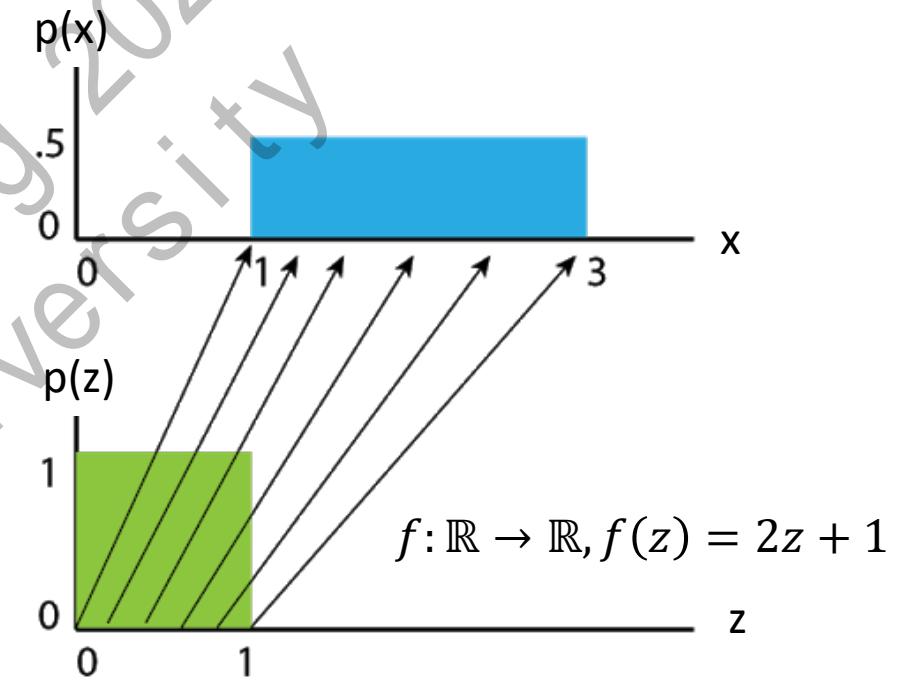
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    - $p(x)dx = p(z)dz$



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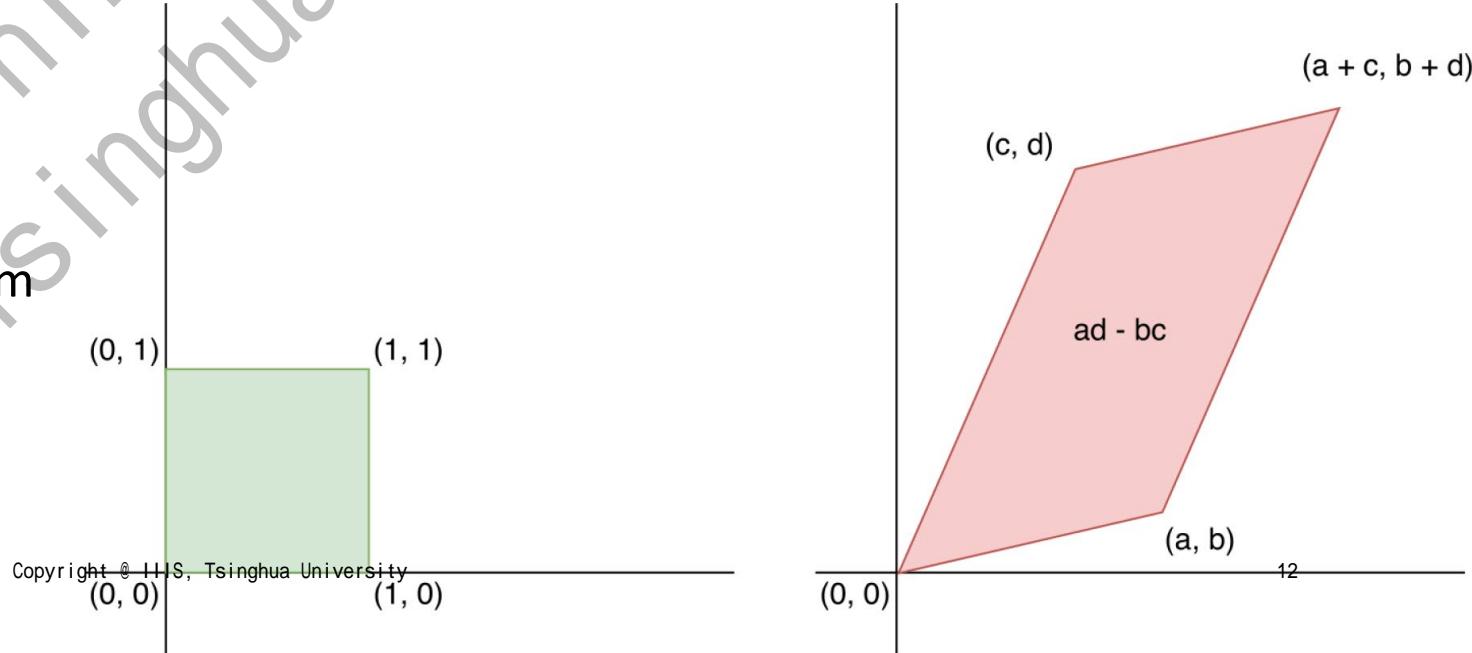


# 2-D Example

- Goal: design  $x = f(z; \theta)$  s.t.
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- Uniform:  $z = [z_1, z_2] \sim \text{unif}([0,1] \times [0,1])$ 
  - Density  $p(z) = 1$
  - $x = Az$ , then  $p(x) = ?$ 
    - $A = \begin{bmatrix} a & b \\ c & d \end{bmatrix}$

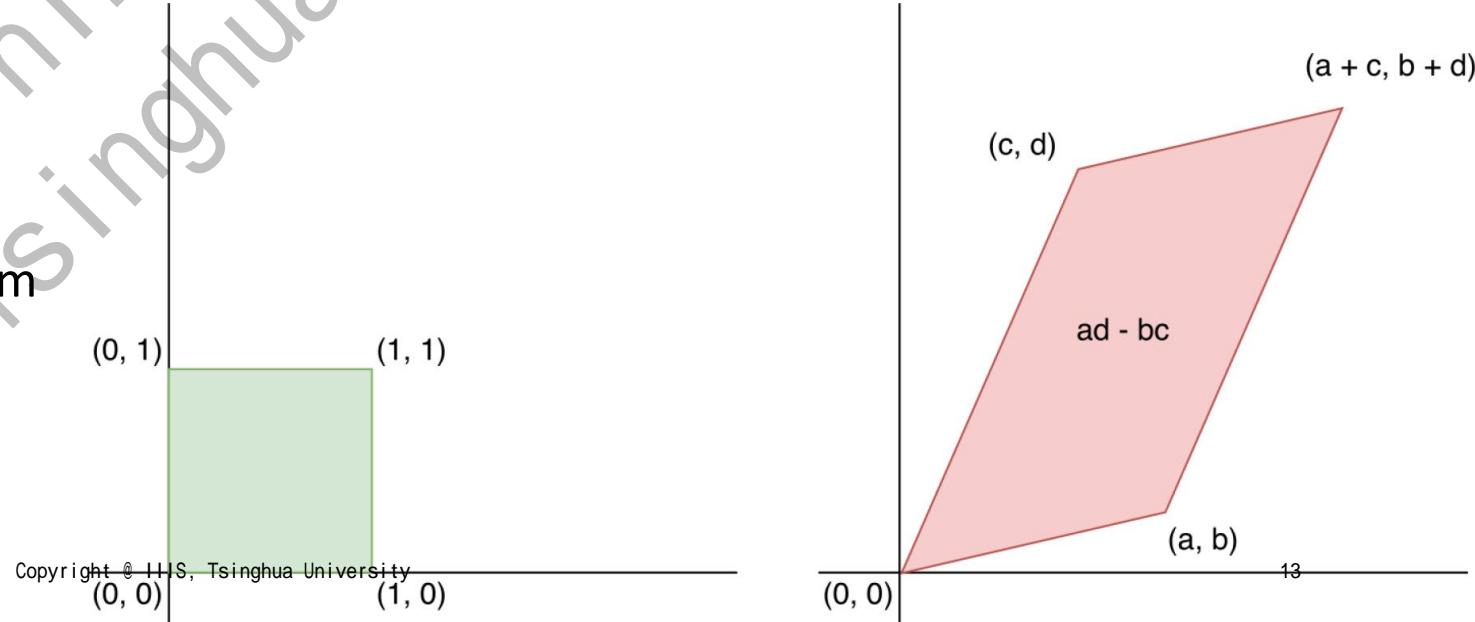
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    - $z$  is mapped to a parallelogram



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  - Density  $p(z) = 1$
  - $x = Az$ , then  $p(x) = 1/S$ 
    - $A = \begin{bmatrix} a & b \\ c & d \end{bmatrix}$
    - $z$  is mapped to a parallelogram
    - $S = |ad - bc|$ , the area



# 2-D Geometry

- The area of the parallelogram is equivalent to the determinant of  $A$

$$\det A = \det \begin{bmatrix} a & b \\ c & d \end{bmatrix} = ad - bc$$

- For any linear transformation  $x = Az + b$ 
  - The following holds (for space of any dimensions)
$$p(x) = |\det A|^{-1} \cdot p(z)$$
  - Remark:  $A$  has a full rank! (bijection)
- More general cases: the change of variable

# Change of Variable

- Suppose  $x = f(z)$  w.r.t. a general non-linear  $f(\cdot)$ 
  - the linearized change in volume is determined by the Jacobian of  $f(\cdot)$

$$\frac{\partial f(z)}{\partial z} = \begin{bmatrix} \frac{\partial f_1(z)}{\partial z_1} & \dots & \frac{\partial f_1(z)}{\partial z_d} \\ \vdots & \ddots & \vdots \\ \frac{\partial f_d(z)}{\partial z_1} & \dots & \frac{\partial f_d(z)}{\partial z_d} \end{bmatrix}$$

- Given a bijection  $f(z): \mathbb{R}^d \rightarrow \mathbb{R}^d$

- $z = f^{-1}(x)$

$$p(x) = p(f^{-1}(x)) \left| \det \left( \frac{\partial f^{-1}(x)}{\partial x} \right) \right| = p(z) \left| \det \left( \frac{\partial f^{-1}(x)}{\partial x} \right) \right|$$

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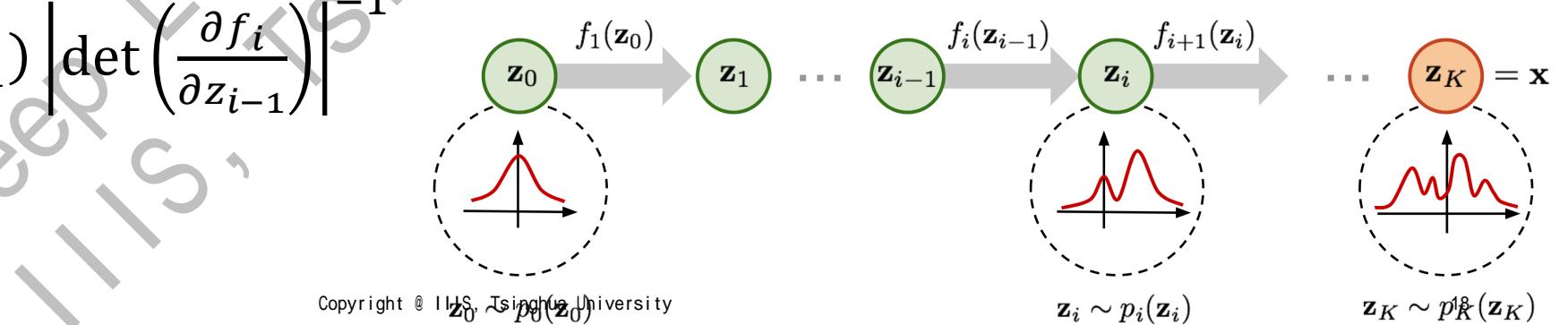
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# Normalizing Flow

- Idea
  - Sample  $z_0$  from an “easy” distribution, i.e., a standard Gaussian
  - Apply  $K$  bijections  $z_i = f_i(z_{i-1})$   $1 \leq i \leq K$
  - The final sample  $x = f_K(z_K)$  has tractable density
- Normalizing Flow

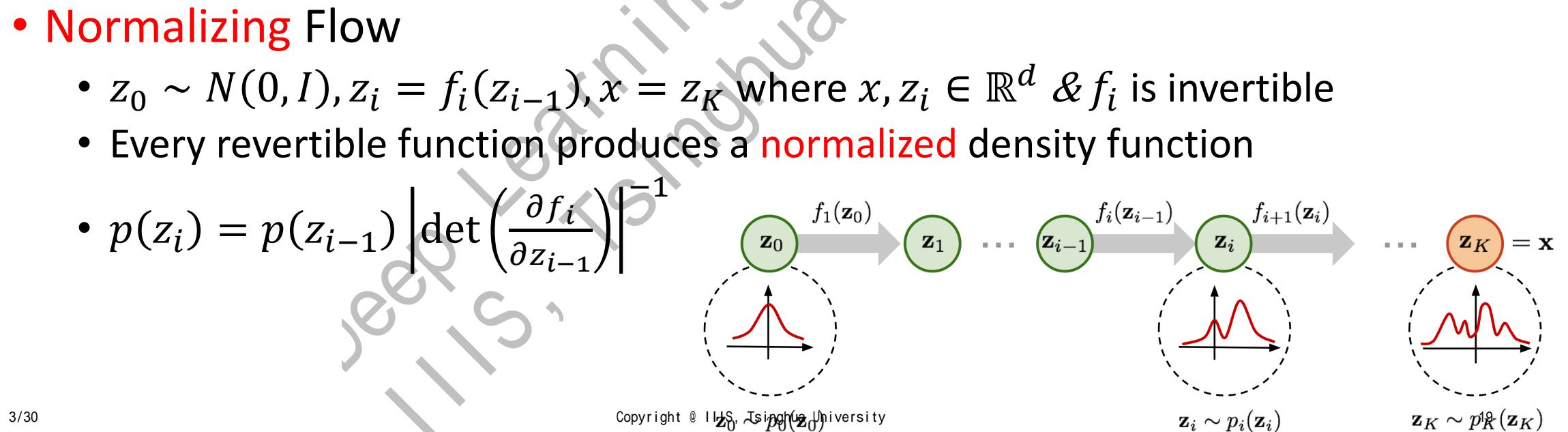
- $z_0 \sim N(0, I)$ ,  $z_i = f_i(z_{i-1})$ ,  $x = z_K$  where  $x, z_i \in \mathbb{R}^d$  &  $f_i$  is invertible
- Every revertible function produces a normalized density function

$$p(z_i) = p(z_{i-1}) \left| \det \left( \frac{\partial f_i}{\partial z_{i-1}} \right) \right|^{-1}$$



# Normalizing Flow

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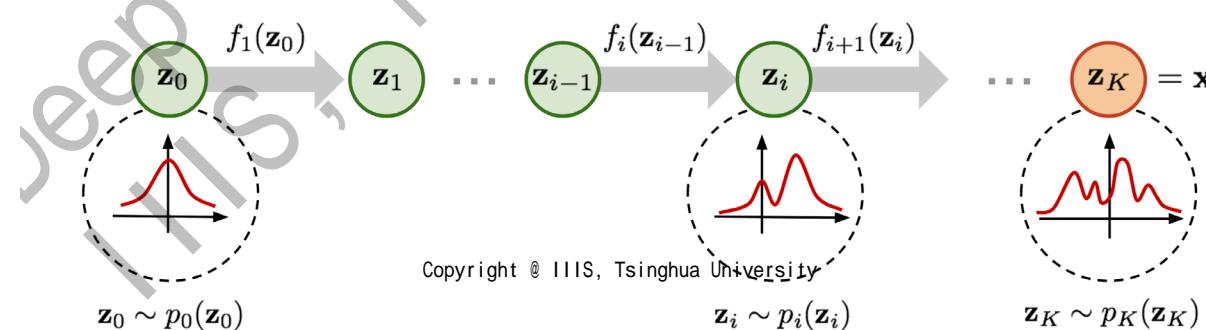


# Normalizing Flow

- Generation is trivial
  - Sample  $z_0$ , then apply the transformations
- Log-Likelihood

$$\log p(x) = \log p(z_{K-1}) - \log \left| \det \left( \frac{\partial f_K}{\partial z_{K-1}} \right) \right|$$

$$\log p(x) = \log p(z_0) - \sum_i \log \left| \det \left( \frac{\partial f_i}{\partial z_{i-1}} \right) \right|$$

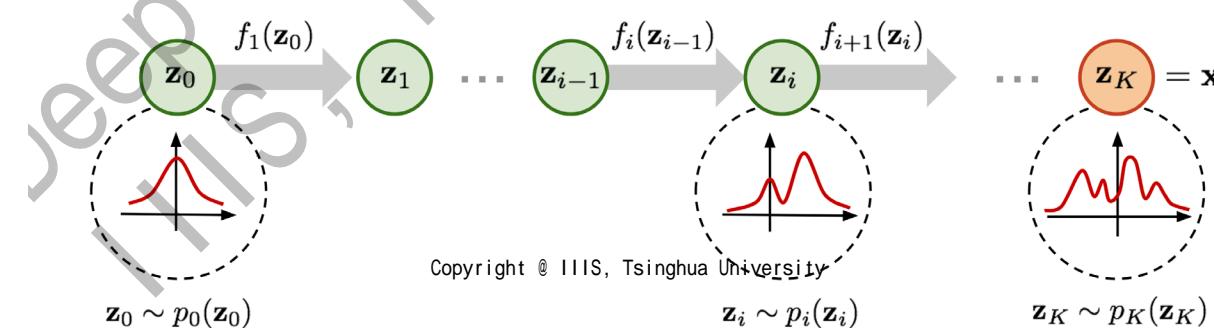


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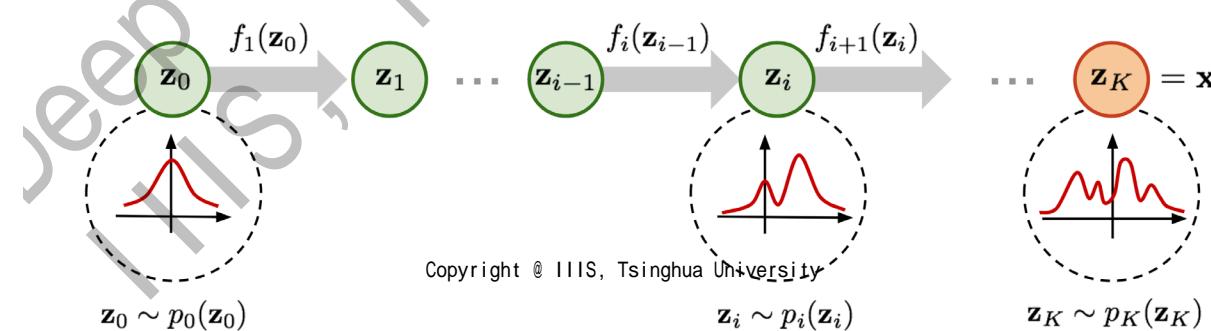
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**$O(d^3)!!!$**



# Normalizing Flow

- Naïve flow model requires extremely expensive computation
  - Determinant of a  $d \times d$  matrix
- Idea
  - Design a good bijection  $f_i(z)$  such that the determinant is easy to compute
- Technical Keys
  - Bjection
    - Randomly constructed matrices are typically full-rank
  - **Structured Jacobian**
    - Desired Jacobian structures for fast determinant computation

# Triangular Jacobian

- Given  $x = (x_1, \dots, x_d) = f(z) = (f_1(z), \dots, f_d(z))$

$$J = \frac{\partial f}{\partial z} = \begin{bmatrix} \frac{\partial f_1}{\partial z_1} & \cdots & \frac{\partial f_1}{\partial z_d} \\ \vdots & \ddots & \vdots \\ \frac{\partial f_d}{\partial z_1} & \cdots & \frac{\partial f_d}{\partial z_d} \end{bmatrix}$$

- Suppose  $x_i = f_i(z)$  only depends on  $z_{\leq i}$ , then

$$\det J = \det \left| \frac{\partial f}{\partial z} \right| = \det \begin{bmatrix} \frac{\partial f_1}{\partial z_1} & \cdots & 0 \\ \vdots & \ddots & \vdots \\ \frac{\partial f_d}{\partial z_1} & \cdots & \frac{\partial f_d}{\partial z_d} \end{bmatrix} = \det \text{diag}(J) = \prod_i \frac{\partial f_i}{\partial z_i}$$

# NICE

- Nonlinear Independent Components Estimation (Dinh et. al, 2014)
  - $z = f(x)$ 
    - Notational convention for MLE learning
    - we partition  $x$  into two disjoint subsets  $x_{1:m}$  and  $x_{m+1:d}$  for any  $1 \leq m \leq d$
    - Forward pass  $x \rightarrow z$  (inference)
      - $z_{1:m} = x_{1:m}$  (identity)
      - $z_{m+1:d} = x_{m+1:d} + \mu_\theta(x_{1:m})$  ( $\mu_\theta$  is a neural network)
    - Backward pass  $z \rightarrow x$  (sampling)
      - $x_{1:m} = z_{1:m}$  (identity)
      - $x_{m+1:d} = z_{m+1:d} - \mu_\theta(z_{1:m})$
  - Volume preserving transformation
    - $\det J = 1$

$$J = \frac{\partial f}{\partial x} = \begin{bmatrix} I_m & 0 \\ \frac{\partial \mu}{\partial x_{1:m}} & I_{d-m} \end{bmatrix}$$

# NICE

- Coupling layers are introduced to ensure all dimensions are covered
  - Reverse (or randomly shuffle) the partition before each transformation layer
- First layer of NICE uses a re-scaling layer
  - $z_i = S_{ii}x_i$
  - Ensure non-unit volume transformation
  - Jacobian of forward pass

$$J = \text{diag}(S)$$
$$\det J = \prod_i S_{ii}$$

# NICE

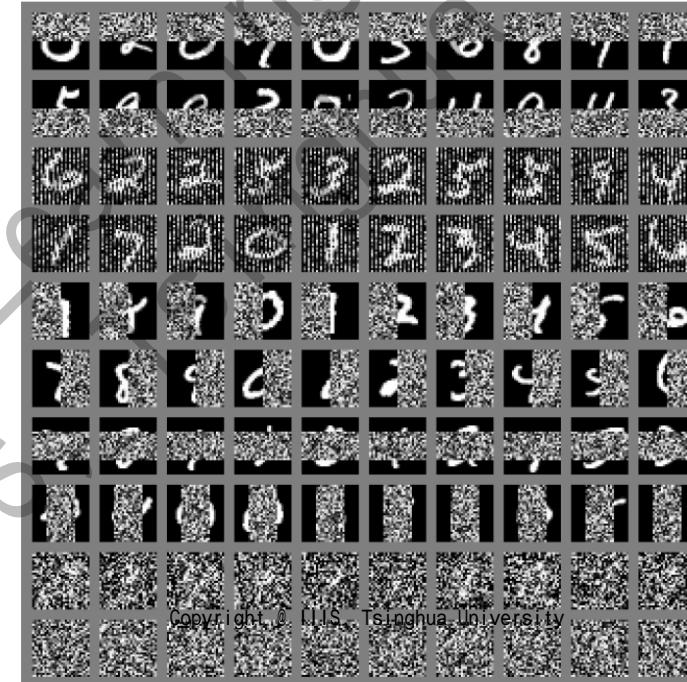
- Generation Results

2	0	0	4	3	5	6	6	8	7
9	5	3	7	4	8	1	6	0	2
5	8	2	3	7	1	3	9	2	6
3	8	7	0	8	6	0	2	3	0
0	3	3	8	2	8	5	0	5	3
9	5	1	3	3	5	9	0	9	7
6	8	5	7	2	4	0	5	6	4
3	5	9	0	8	7	4	7	3	1
3	6	2	0	0	9	3	3	4	2
9	2	4	8	3	6	9	8	7	6



# NICE

- Inpainting
  - $x = (x_v, x_h)$
  - We have tractable likelihood function  $p(x_v, x_h)$ !
    - Gradient ascent (stochastic gradient MCMC if you want samples)



# Real-NVP

- NICE: most layers maintain an *unchanged* volume
- Non-volume preserving extension of NICE (Dinh et al, 2016)
  - Two partitions over  $z$ :  $z_{1:m}$  and  $z_{m+1:d}$  for any  $1 \leq m \leq d$
  - Forward pass  $x \rightarrow z$  (inference)
    - $z_{1:m} = x_{1:m}$  (identity)
    - $z_{m+1:d} = x_{m+1:d} \cdot \exp(\alpha_\theta(x_{1:m})) + \mu_\theta(x_{1:m})$  ( $\mu_\theta$  &  $\alpha_\theta$  are neural networks)
  - Backward pass  $z \rightarrow x$  (sampling)
    - $x_{1:m} = z_{1:m}$  (identity)
    - $x_{m+1:d} = (z_{m+1:d} - \mu_\theta(z_{1:m})) \cdot \exp(-\alpha_\theta(x_{1:m}))$
  - Non-volume preserving transformation

$$\det J = \prod_{i=m+1}^d \exp(\alpha_\theta(x_{1:m})_i)$$

# Real-NVP

- Generation Results
  - Left: training data; Right: generated samples



# Real-NVP

- Explore the latent space
  - 4 samples selected:  $z^0, z^1, z^2, z^3$ , two interpolation parameters  $\alpha, \beta$
  - $z = \cos(\alpha) (\cos(\beta)z^1 + \sin(\beta)z^2) + \sin(\alpha)(\cos(\beta)z^3 + \sin(\beta)z^4)$



# Real-NVP

- Fun Fact

Accepted as a workshop contribution at ICLR 2015

Published as a conference paper at ICLR 2017

## NICE: NON-LINEAR INDEPENDENT COMPONENTS ESTIMATION

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## DENSITY ESTIMATION USING REAL NVP

**Laurent Dinh\***

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University of Montreal  
Montreal, QC H3T1J4

**Jascha Sohl-Dickstein**

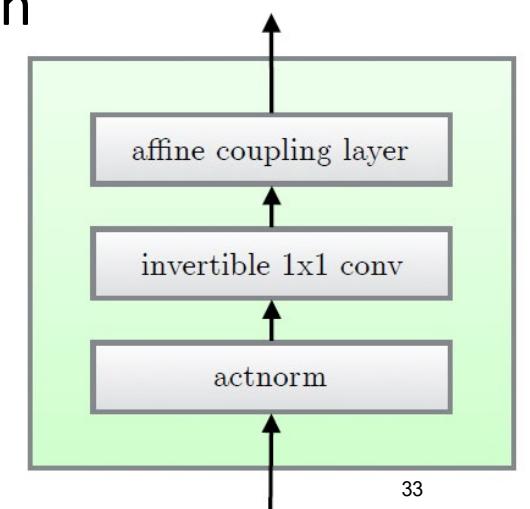
Google Brain

**Samy Bengio**

Google Brain

# GLOW

- Limited expressiveness of previous coupling layers
  - But a general non-linear transformation can be too expensive...
- Generative Flow with Invertible 1x1 Convolutions (Kingma et al. 2018)
  - Input:  $x = h \times w \times c$  (height, width, channel) (assume  $c$  is small)
  - Key idea: introduce 1x1 convolutions when channel size is small
  - 1x1 conv: a linear transformation for each feature map location
    - Forward mapping:  $z_{ij} = Wx_{ij} + b$
    - Inverse mapping: simply compute the inverse matrix of  $W$
  - Computation  $O(c^3)$ 
    - $\log|\det J| = h \cdot w \cdot \log |\det W|$
  - Also use normalization layer to stabilize training
    - Architecture details can be found in the paper



# GLOW

- Generation Results



# Normalizing Flow: Summary

- Key Ideas
  - Generation by iteratively transforming a simple distribution
  - Invertible transformation for tractable likelihood
    - Enable straightforward MLE learning
  - Design principle
    - Apply non-linear transformations with easy-to-compute Jacobian determinants
- Pros & Cons
  - Easy sampling & training via deterministic transformation from a simple distribution
  - Most restricted network structure (trade expressiveness for tractability)
    - Architecture, dimensionality, etc.
    - Most suitable for the use cases where tractability is a must

# Today's Topic

- Normalizing Flow
  - A generative model class that has the best sampling and likelihood properties
- Score-based generative model
  - A different framework to tackle general energy-based models

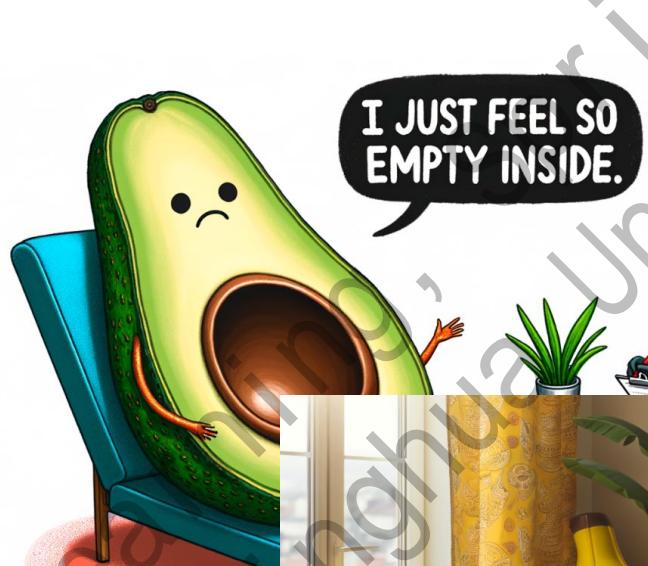
# Today's Topic

- Normalizing Flow
  - A generative model class that has the best sampling and likelihood properties
- Score-based generative model
  - A different framework to tackle general energy-based models
  - The model class that has the best generation quality
  - It is also called the *diffusion model*

# Why Diffusion Model?



Dall-E 3



# Why Diffusion Model?



Stable Diffusion Model (SD3)

# Why Diffusion



Midjourney



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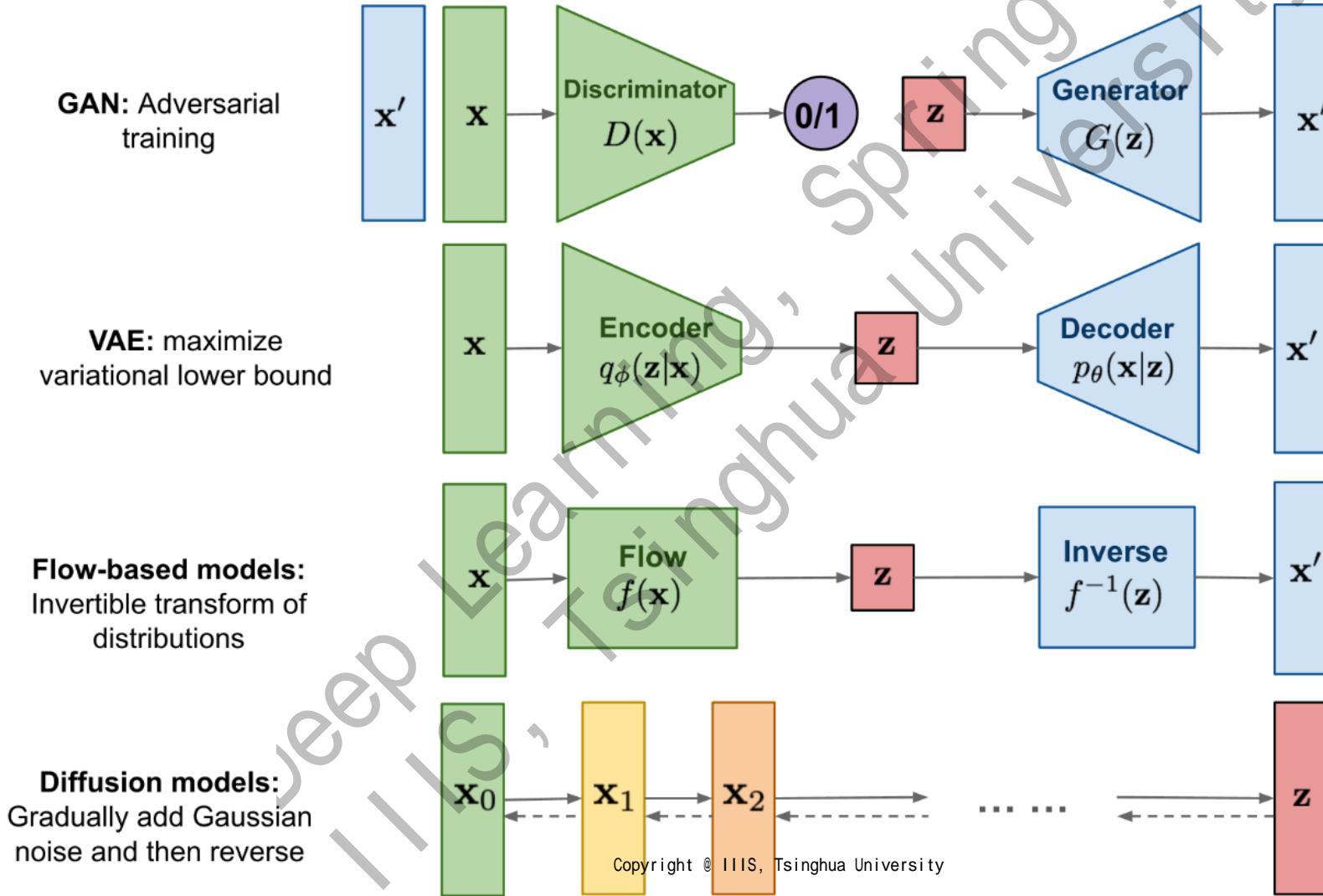
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# Why Diffusion Model?

- AI-generated photo wins the award
  - Jason Allen's A.I.-generated work, "Théâtre D'opéra Spatial," took first place in the digital category at the Colorado State Fair.
- The trend of AIGC
  - AI Generated Content
- Foundation: Diffusion Model



# What is Diffusion Model



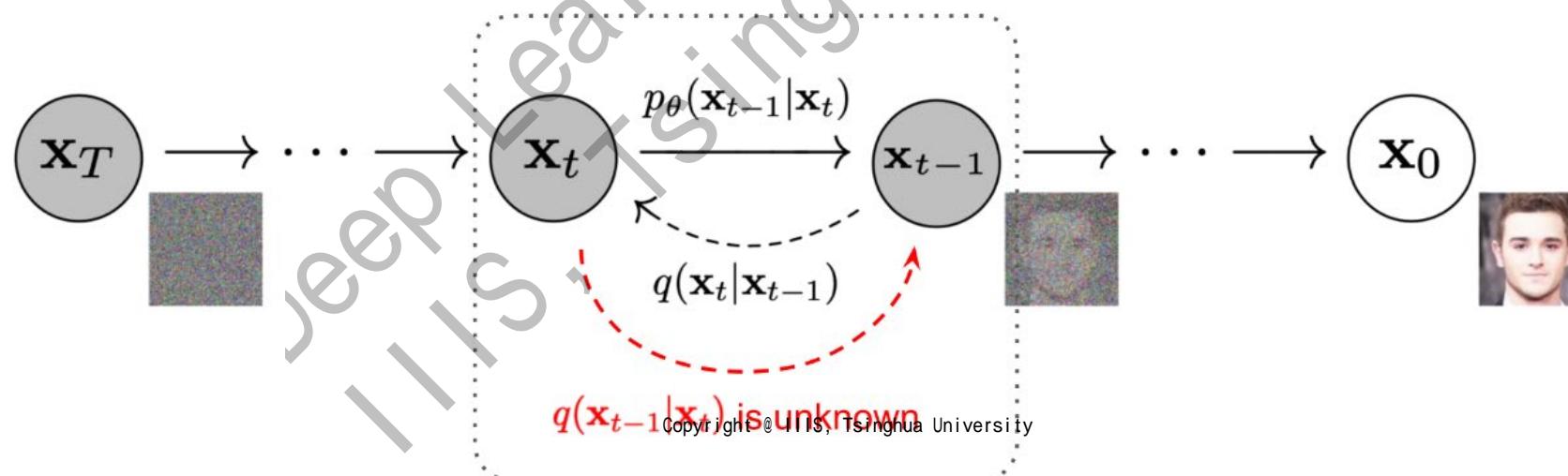
# What is Diffusion Model

- Formal Definition

- $x_0 \sim q_{data}(x)$
- Forward diffusion process: continuously adding Gaussian noise to data

$$q(\mathbf{x}_t | \mathbf{x}_{t-1}) = \mathcal{N}(\mathbf{x}_t; \sqrt{1 - \beta_t} \mathbf{x}_{t-1}, \beta_t \mathbf{I}) \quad q(\mathbf{x}_{1:T} | \mathbf{x}_0) = \prod_{t=1}^T q(\mathbf{x}_t | \mathbf{x}_{t-1})$$

- Sampling process: gradually recover the data from isomorphic Gaussian noise



# Diffusion Model

- Milestone works
  - The story
    - <https://www.quantamagazine.org/the-physics-principle-that-inspired-modern-ai-art-20230105/>
  - Deep Unsupervised Learning using Nonequilibrium Thermodynamics (ICML 2015)
    - The original diffusion model
  - Generative Modeling by Estimating Gradients of the Data Distribution (Yang Song, et al., NIPS 2019)
    - Score-based model, foundation of modern diffusion model
  - Denoising Diffusion Probabilistic Models (Jonathan Ho, et al., NIPS 2020)
    - DDPM: the first working diffusion model

# Diffusion Model

- Milestone works
  - The story
    - <https://www.quantamagazine.org/the-physics-principle-that-inspired-modern-ai-art-20230105/>
  - Deep Unsupervised Learning using Nonequilibrium Thermodynamics (ICML 2015)
    - The original diffusion model
  - Generative Modeling by Estimating Gradients of the Data Distribution (Yang Song, et al., NIPS 2019)
    - Score-based model, foundation of modern diffusion model
  - **Denoising Diffusion Probabilistic Models (Jonathan Ho, et al., NIPS 2020)**
    - DDPM: the first working diffusion model

# Diffusion Model

- Denoising Diffusion Probabilistic Models (Jonathan Ho, et al., NIPS 2020)
  - Learning  $\epsilon_\theta(x, t)$ ;  $\bar{\alpha}_t = \prod_{i=1}^t \alpha_i$
  - Original diffusion model loss function from ICML15 (your homework 😊)
    - $L_t = E_{x_0, \epsilon} \left[ \frac{(1-\alpha_t)^2}{2\alpha_t(1-\bar{\alpha}_t)|\Sigma_\theta|_2^2} |\epsilon_t - \epsilon_\theta(\sqrt{\bar{\alpha}_t}x_0 + \sqrt{1-\bar{\alpha}_t}\epsilon_t, t)|^2 \right]$
  - DDPM simplified objective:  $T = 1000, \alpha_t = 1 - \beta_t, \beta_t \sim [10^{-4}, 0.02]$

---

**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  $\mathbf{x}_0$ 
  
```

**Why does it work?**

# Diffusion Model

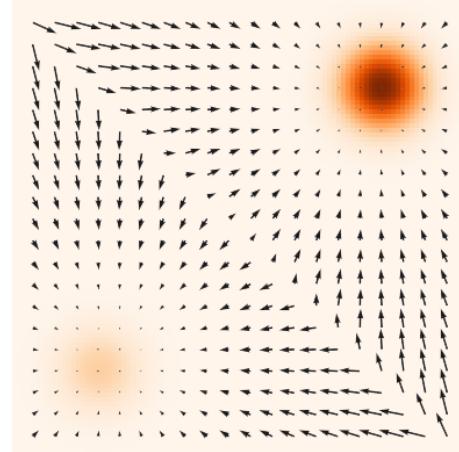
- Milestone works
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# Diffusion Model

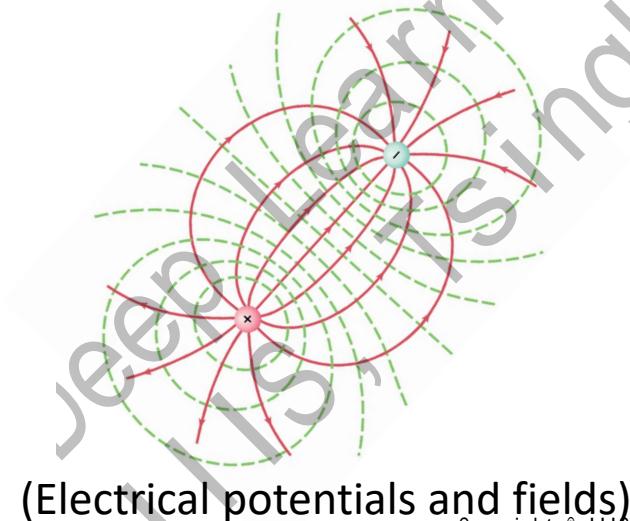
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    - **Score-based model, foundation of modern diffusion model**
  - Denoising Diffusion Probabilistic Models (Jonathan Ho, et al., NIPS 2020)
    - DDPM: the first working diffusion model
    - **A simplified training objective directly connected to score-based model**

# Score-Based Model

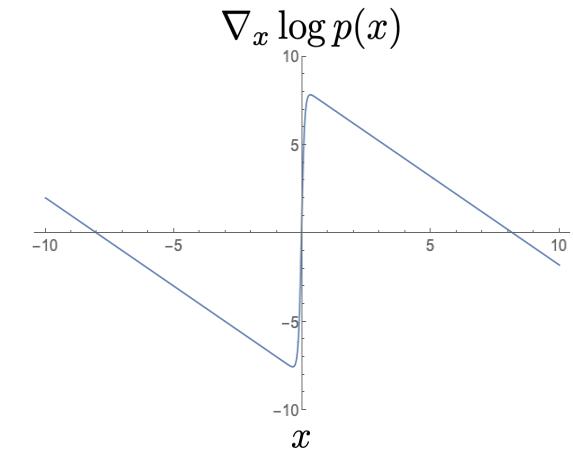
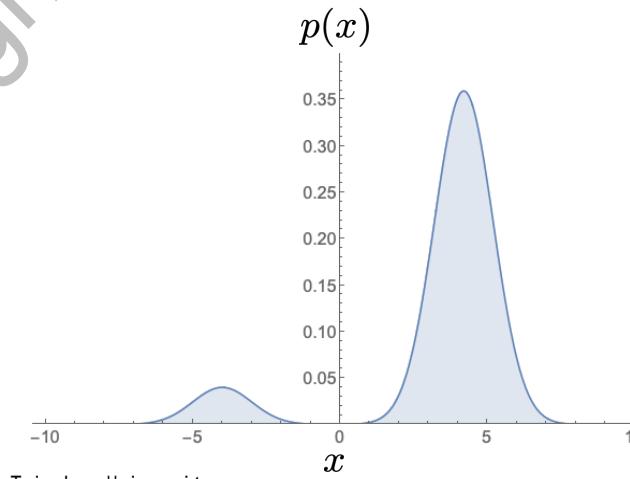
- How to represent a distribution  $p(x)$ ?
  - When the pdf is differentiable, we can compute the gradient of a probability density.
  - **Score function**:  $s(x) = \nabla_x \log p(x)$



(pdf and score)



(Electrical potentials and fields)



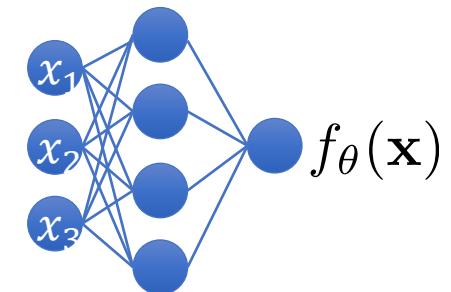
# Score-Based Model

- Energy-based model (recap)

- Energy function  $f_\theta(\mathbf{x})$
- Partition function  $Z(\theta)$

$$f_\theta(\mathbf{x}) \in \mathbb{R}$$

$$p_\theta(\mathbf{x}) = \frac{e^{f_\theta(\mathbf{x})}}{Z(\theta)}$$



- Learning EBMs for  $p_{data}$

- MLE with Contrastive Divergence for  $\mathbf{x}_{\text{train}} \sim p_{data}$

$$\max_{\theta} f_\theta(\mathbf{x}_{\text{train}}) - \log Z(\theta)$$

$$\nabla_{\theta} f_\theta(\mathbf{x}_{\text{train}}) - \nabla_{\theta} \log Z(\theta) \approx \nabla_{\theta} f_\theta(\mathbf{x}_{\text{train}}) - \nabla_{\theta} f_\theta(\mathbf{x}_{\text{sample}})$$

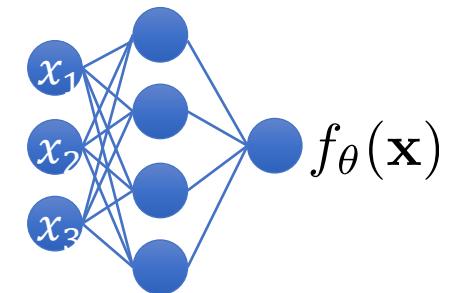
- Monte-Carlo sampling for negative samples  $\mathbf{x}_{\text{sample}} \sim p_\theta(\mathbf{x})$

# Score-Based Model

- Energy-based model (recap)
  - Energy function  $f_\theta(x)$
  - Partition function  $Z(\theta)$
- Learning EBMs for  $p_{data}$ 
  - An alternative objective: Score matching
    - fisher divergence  $F(p||q) = \frac{1}{2} E_{x \sim p} [ \|\nabla_x p(x) - \nabla_x q(x) \|_2^2 ]$
    - Score matching by minimizing fisher divergence

$$f_\theta(\mathbf{x}) \in \mathbb{R}$$

$$p_\theta(\mathbf{x}) = \frac{e^{f_\theta(\mathbf{x})}}{Z(\theta)}$$



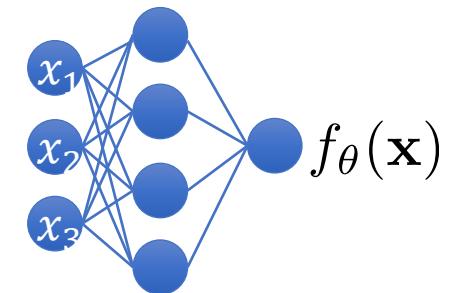
$$\begin{aligned} & \frac{1}{2} E_{x \sim p_{data}} [ \|\nabla_x \log p_{data}(x) - \nabla_x \log p_\theta(x) \|_2^2 ] \\ &= E_{x \sim p_{data}} \left[ \frac{1}{2} \|\nabla_x \log p_\theta(x)\|_2^2 + \text{tr}(\nabla_x^2 \log p_\theta(x)) \right] + \text{Const} \end{aligned}$$

# Score-Based Model

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$$\begin{aligned} & \frac{1}{2} E_{x \sim p_{data}} [ \|\nabla_x \log p_{data}(x) - \nabla_x \log p_\theta(x) \|_2^2 ] \\ &= E_{x \sim p_{data}} \left[ \frac{1}{2} \|\nabla_x \log p_\theta(x)\|_2^2 + \text{tr}(\nabla_x^2 \log p_\theta(x)) \right] + \text{Const} \end{aligned}$$

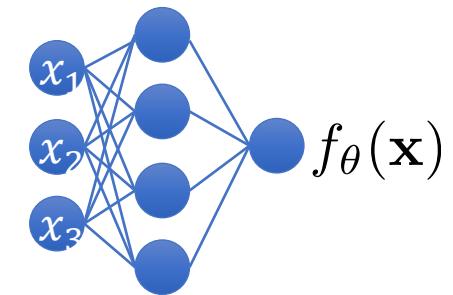
# Score-Based Model

- Energy-based model (recap)

- Energy function  $f_\theta(\mathbf{x})$
- Partition function  $Z(\theta)$

$$f_\theta(\mathbf{x}) \in \mathbb{R}$$

$$p_\theta(\mathbf{x}) = \frac{e^{f_\theta(\mathbf{x})}}{Z(\theta)}$$



- Learning EBMs for  $p_{data}$

- An alternative objective: Score matching

- fisher divergence  $F(p||q) = \frac{1}{2} E_{x \sim p} [ \|\nabla_x p(x) - \nabla_x q(x) \|_2^2 ]$

- Score matching by minimizing fisher divergence

$$\frac{1}{2} E_{x \sim p_{data}} [ \|\nabla_x \log p_{data}(x) - \nabla_x \log p_\theta(x) \|_2^2 ]$$

$$= E_{x \sim p_{data}} \left[ \frac{1}{2} \|\nabla_x f_\theta(x)\|_2^2 + \text{tr}(\nabla_x^2 f_\theta(x)) \right] + \text{Const}$$

**No Partition function any more**

# Score-Based Model

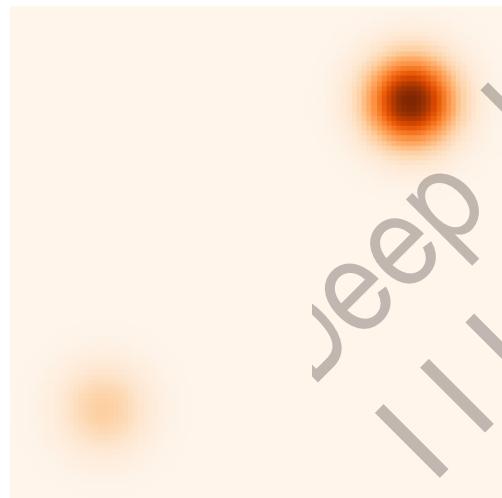
- Score-based model is beyond EBM

- $s_\theta(x) \approx \nabla_x \log p_{\text{data}}(x)$

- Learning?

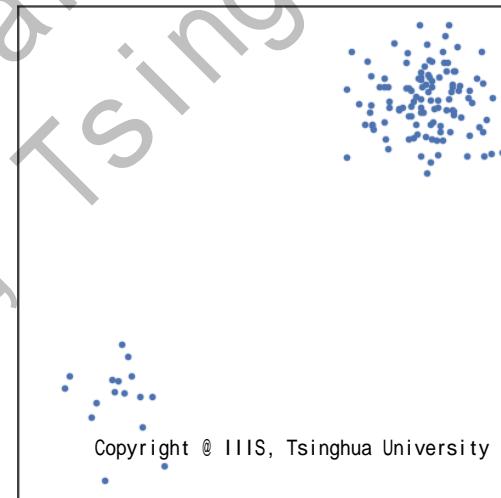
Probability density

$$p_{\text{data}}(\mathbf{x})$$



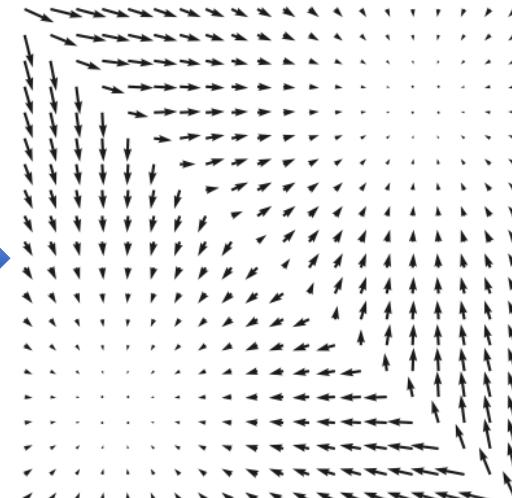
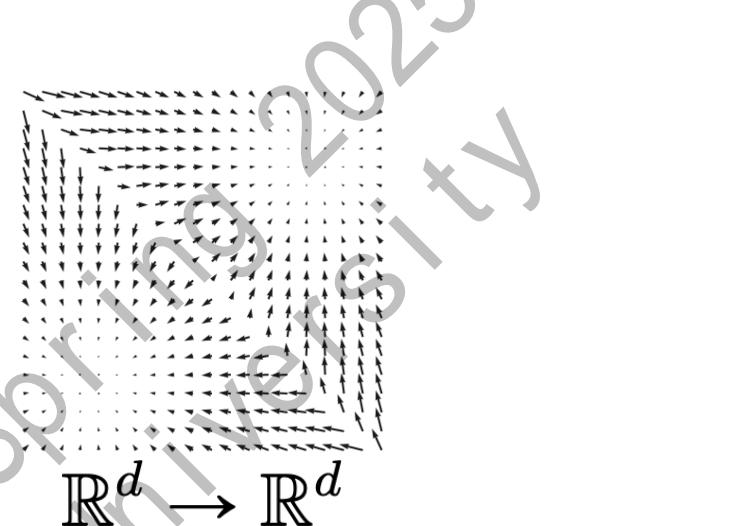
i.i.d. samples

$$\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n$$



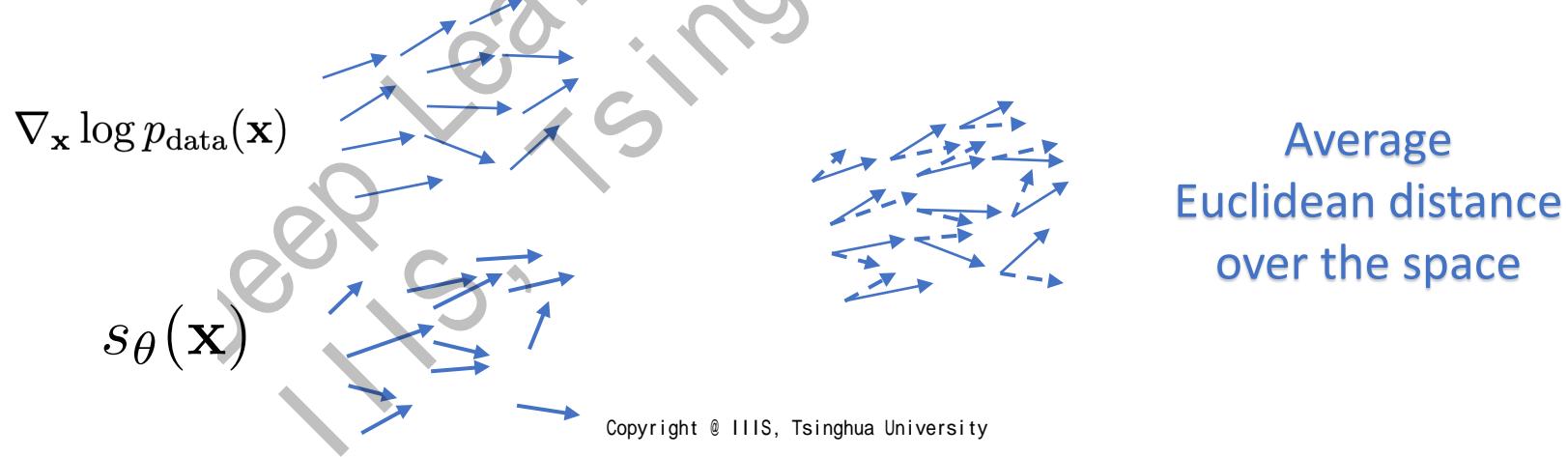
Score function

$$s_\theta(\mathbf{x}) \approx \nabla_{\mathbf{x}} \log p_{\text{data}}(\mathbf{x})$$



# Score-Based Model

- Score estimation formulation
  - Given: i.i.d. samples  $\{x_1, x_2, \dots, x_n\} \sim p_{data}(x)$
  - Task: Estimating the score  $\nabla_x \log p_{data}(x)$
  - Score model: A learnable vector-valued function  $s_\theta(x): \mathbb{R}^d \rightarrow \mathbb{R}^d$
  - Goal:  $s_\theta(x) \approx \nabla_x p_{data}(x)$



# Score-Based Model

- **Objective:** Average Euclidean distance over the whole space.

$$\frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}} [\|\nabla_{\mathbf{x}} \log p_{\text{data}}(\mathbf{x}) - \mathbf{s}_{\theta}(\mathbf{x})\|_2^2]$$

(Fisher divergence)

- **Score matching:**

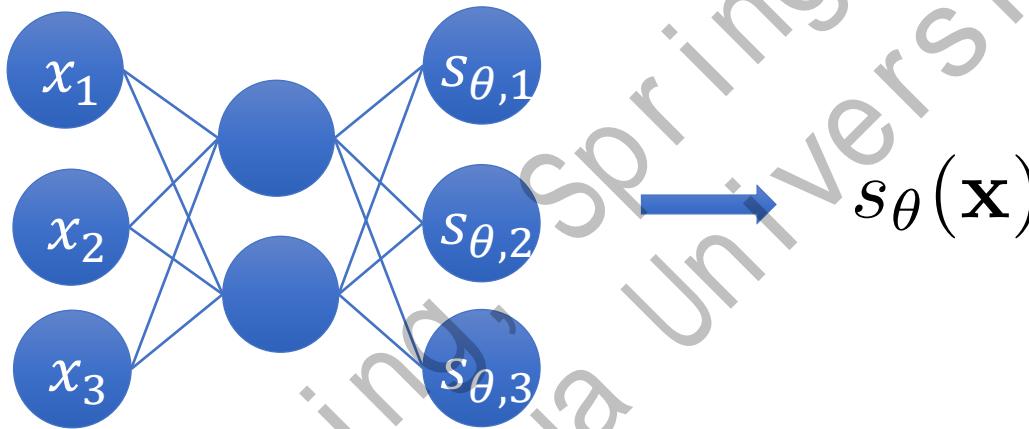
$$E_{\mathbf{x} \sim p_{\text{data}}} \left[ \frac{1}{2} \|\mathbf{s}_{\theta}(\mathbf{x})\|_2^2 + \text{tr} \left( \underbrace{\nabla_{\mathbf{x}} \mathbf{s}_{\theta}(\mathbf{x})}_{\text{Jacobian of } \mathbf{s}_{\theta}(\mathbf{x})} \right) \right]$$

- **Requirements:**

- The score model must be efficient to evaluate.
- **How to have a proper model for the score function?**

## Score-Based Model

- Deep neural networks as more expressive score models



Score Matching  
is not Scalable due  
to the Jacobian!

- Compute  $\|s_\theta(\mathbf{x})\|_2^2$  and  $\text{tr}(\nabla_{\mathbf{x}} s_\theta(\mathbf{x}))$

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# $O(d)$ Backprops!

# Denoising score matching

- Denoising score matching (Vincent 2011): matching the score of a noise-perturbed distribution



$\mathbf{x}$



$\tilde{\mathbf{x}}$

$$p_{\text{data}}(\mathbf{x})$$

$$q_{\sigma}(\tilde{\mathbf{x}} \mid \mathbf{x}) = \text{const.} + \frac{1}{2} E_{\tilde{\mathbf{x}} \sim q_{\sigma}} [\|\mathbf{s}_{\theta}(\tilde{\mathbf{x}})\|_2^2]$$

$$q_{\sigma}(\tilde{\mathbf{x}}) = \text{const.} + \frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}(\mathbf{x}), \tilde{\mathbf{x}} \sim q_{\sigma}(\tilde{\mathbf{x}} \mid \mathbf{x})} [\|\mathbf{s}_{\theta}(\tilde{\mathbf{x}}) - \nabla_{\tilde{\mathbf{x}}} \log q_{\sigma}(\tilde{\mathbf{x}} \mid \mathbf{x})\|_2^2]$$

$$= \text{const.} + \frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}(\mathbf{x}), \tilde{\mathbf{x}} \sim q_{\sigma}(\tilde{\mathbf{x}} \mid \mathbf{x})} [\|\mathbf{s}_{\theta}(\tilde{\mathbf{x}}) - \nabla_{\tilde{\mathbf{x}}} \log q_{\sigma}(\tilde{\mathbf{x}} \mid \mathbf{x})\|_2^2]$$

$$= \text{const.} + \frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}(\mathbf{x}), \tilde{\mathbf{x}} \sim q_{\sigma}(\tilde{\mathbf{x}} \mid \mathbf{x})} [\|\mathbf{s}_{\theta}(\tilde{\mathbf{x}}) - \nabla_{\tilde{\mathbf{x}}} \log q_{\sigma}(\tilde{\mathbf{x}} \mid \mathbf{x})\|_2^2] + \text{const.}$$

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$$\frac{1}{2} E_{\tilde{\mathbf{x}} \sim q_{\sigma}} [\|\nabla_{\tilde{\mathbf{x}}} \log q_{\sigma}(\tilde{\mathbf{x}}) - \mathbf{s}_{\theta}(\tilde{\mathbf{x}})\|_2^2]$$

$$= \text{const.} + \frac{1}{2} E_{\tilde{\mathbf{x}} \sim q_{\sigma}} [\|\mathbf{s}_{\theta}(\tilde{\mathbf{x}})\|_2^2] - \int q_{\sigma}(\tilde{\mathbf{x}}) \nabla_{\tilde{\mathbf{x}}} \log q_{\sigma}(\tilde{\mathbf{x}})^T \mathbf{s}_{\theta}(\tilde{\mathbf{x}}) d\tilde{\mathbf{x}}$$

Your homework ☺

# Denoising score matching

- Estimate the score of a noise-perturbed distribution

$$\begin{aligned} & \frac{1}{2} E_{\tilde{\mathbf{x}} \sim p_{\text{data}}} [\| \mathbf{s}_\theta(\tilde{\mathbf{x}}) - \nabla_{\tilde{\mathbf{x}}} \log q_\sigma(\tilde{\mathbf{x}}) \|_2^2] \\ &= \frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}(\mathbf{x}), \tilde{\mathbf{x}} \sim q_\sigma(\tilde{\mathbf{x}} | \mathbf{x})} [\| \mathbf{s}_\theta(\tilde{\mathbf{x}}) - \nabla_{\tilde{\mathbf{x}}} \log q_\sigma(\tilde{\mathbf{x}} | \mathbf{x}) \|_2^2] + \text{const.} \end{aligned}$$

- $\nabla_{\tilde{\mathbf{x}}} \log q_\sigma(\tilde{\mathbf{x}} | \mathbf{x})$  is easy to compute
  - $q_\sigma(\tilde{\mathbf{x}} | \mathbf{x}) = \mathcal{N}(\tilde{\mathbf{x}} | \mathbf{x}, \sigma^2 \mathbf{I})$
  - $\nabla_{\tilde{\mathbf{x}}} \log q_\sigma(\tilde{\mathbf{x}} | \mathbf{x}) = -\frac{\tilde{\mathbf{x}} - \mathbf{x}}{\sigma^2}$
- **Pros:** efficient to optimize even for very high dimensional data, and useful for optimal denoising.
- **Con:** cannot estimate the score of clean data (noise-free)

# Denoising score matching

- Sample a minibatch of datapoints  $\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\} \sim p_{\text{data}}(\mathbf{x})$
- Sample a minibatch of perturbed datapoints  $\{\tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \dots, \tilde{\mathbf{x}}_n\} \sim q_\sigma(\tilde{\mathbf{x}})$   

$$\tilde{\mathbf{x}}_i \sim q_\sigma(\tilde{\mathbf{x}}_i \mid \mathbf{x}_i)$$
- Estimate the denoising score matching loss with empirical means

$$\frac{1}{2n} \sum_{i=1}^n [\|s_\theta(\tilde{\mathbf{x}}_i) - \nabla_{\tilde{\mathbf{x}}} \log q_\sigma(\tilde{\mathbf{x}}_i \mid \mathbf{x}_i)\|_2^2]$$

- If Gaussian perturbation

$$\frac{1}{2n} \sum_{i=1}^n \left[ \|s_\theta(\tilde{\mathbf{x}}_i) + \frac{\tilde{\mathbf{x}}_i - \mathbf{x}_i}{\sigma^2}\|_2^2 \right]$$

- Stochastic gradient descent
- Need to choose a very small  $\sigma$ !

# Pitfall of denoising score matching

- The loss variance will increase drastically as  $\sigma \rightarrow 0$ !
- Denoising score matching loss for Gaussian perturbations

$$\begin{aligned}
 & \frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}} E_{\tilde{\mathbf{x}} \sim q_{\sigma}(\tilde{\mathbf{x}} | \mathbf{x})} \left[ \left\| s_{\theta}(\tilde{\mathbf{x}}) + \frac{\tilde{\mathbf{x}} - \mathbf{x}}{\sigma^2} \right\|_2^2 \right] \\
 &= \frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}} E_{\mathbf{z} \sim \mathcal{N}(0, \mathbf{I})} \left[ \left\| s_{\theta}(\mathbf{x} + \sigma \mathbf{z}) + \frac{\mathbf{z}}{\sigma} \right\|_2^2 \right] \quad (\text{reparameterization trick}) \\
 &= \frac{1}{2} E_{\mathbf{x} \sim p_{\text{data}}} E_{\mathbf{z} \sim \mathcal{N}(0, \mathbf{I})} \left[ \|s_{\theta}(\mathbf{x} + \sigma \mathbf{z})\|_2^2 + 2s_{\theta}(\mathbf{x} + \sigma \mathbf{z})^T \frac{\mathbf{z}}{\sigma} + \frac{\|\mathbf{z}\|_2^2}{\sigma^2} \right]
 \end{aligned}$$

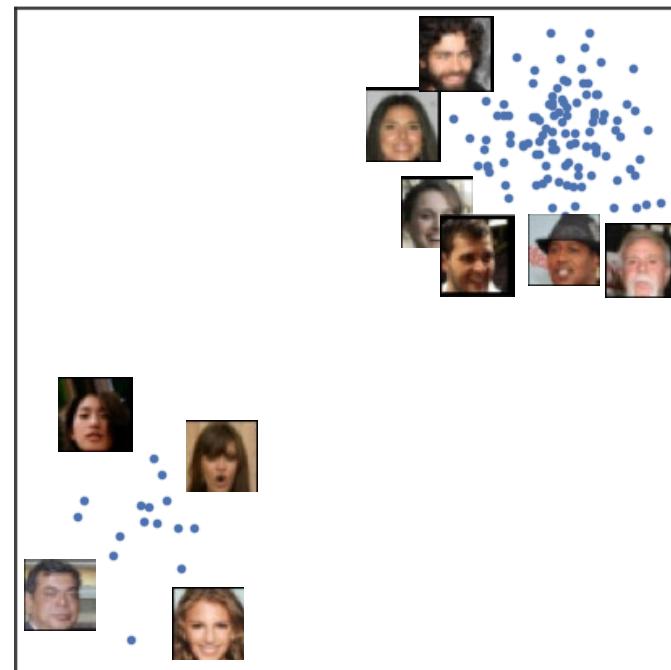
- If we choose very small  $\sigma \rightarrow 0$

$$\text{Var} \left( \frac{\mathbf{z}}{\sigma} \right) \rightarrow \infty$$

$$\text{Var} \left( \frac{\|\mathbf{z}\|_2^2}{\sigma^2} \right) \rightarrow \infty$$

**We need to tune  $\sigma$  carefully!**

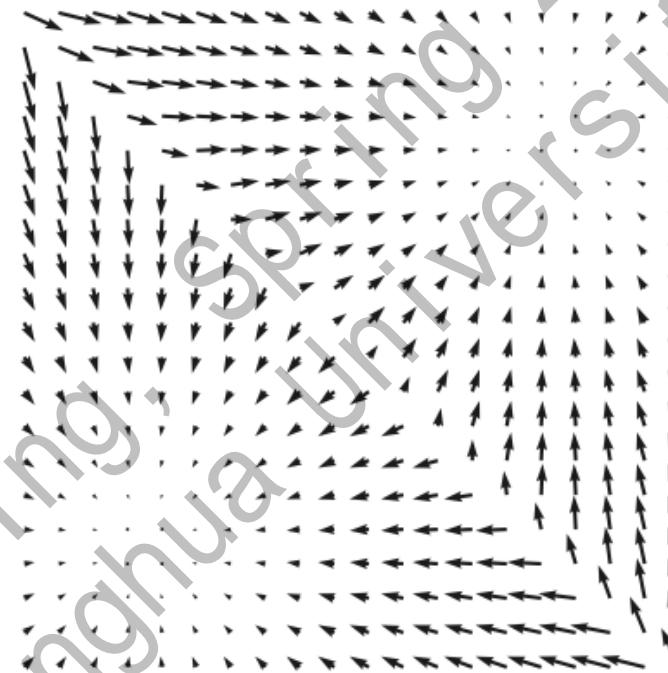
# Score-based generative modeling



Data samples

$$\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\} \sim p_{\text{data}}(\mathbf{x})$$

Score  
Matching



Scores

$$s_{\theta}(\mathbf{x}) \approx \nabla_{\mathbf{x}} \log p_{\text{data}}(\mathbf{x})$$

?

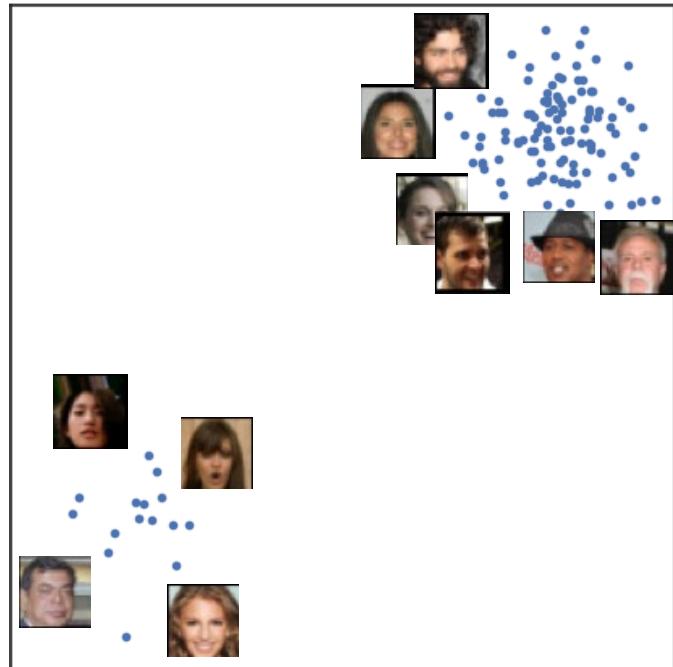


New samples

# Langevin dynamics sampling (Recap)

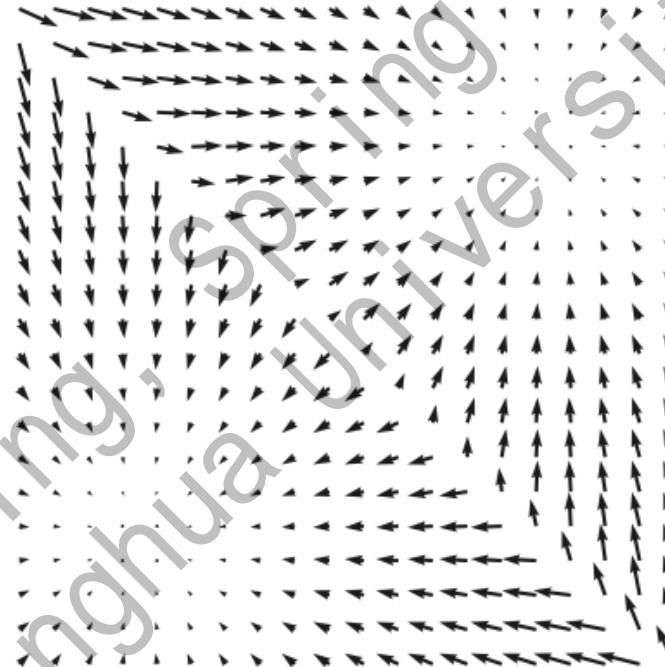
- Sample from  $p(\mathbf{x})$  using only the score  $\nabla_{\mathbf{x}} \log p(\mathbf{x})$
- Initialize  $\mathbf{x}^0 \sim \pi(\mathbf{x})$
- Repeat for  $t \leftarrow 1, 2, \dots, T$   
 $\mathbf{z}^t \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$   
$$\mathbf{x}^t \leftarrow \mathbf{x}^{t-1} + \frac{\epsilon}{2} \nabla_{\mathbf{x}} \log p(\mathbf{x}^{t-1}) + \sqrt{\epsilon} \mathbf{z}^t$$
- If  $\epsilon \rightarrow 0$  and  $T \rightarrow \infty$ , we are guaranteed to have  $\mathbf{x}^T \sim p(\mathbf{x})$
- Langevin dynamics + score estimation  $s_{\theta}(\mathbf{x}) \approx \nabla_{\mathbf{x}} \log p(\mathbf{x})$

# Score-based generative modeling

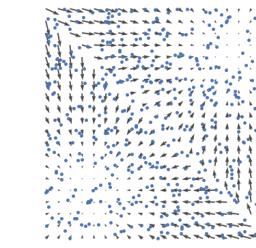


$$\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\} \sim p_{\text{data}}(\mathbf{x})$$

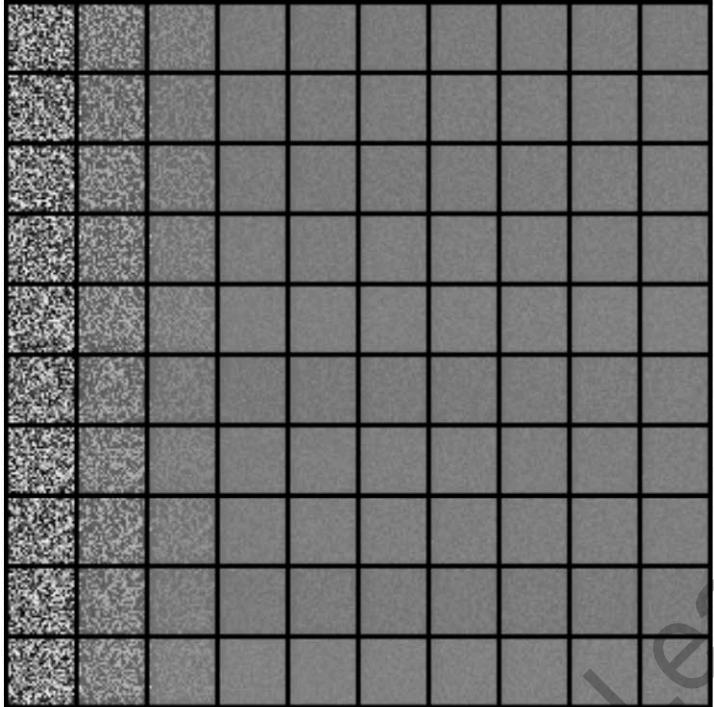
score  
matching



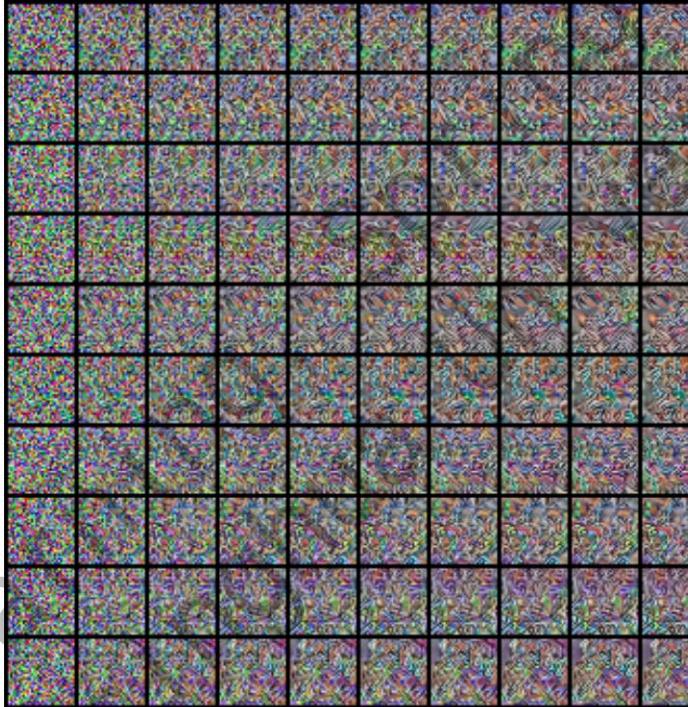
$$s_\theta(\mathbf{x}) \approx \nabla_{\mathbf{x}} \log p_{\text{data}}(\mathbf{x})$$



# Score-based generative modeling: empirical results



(a) MNIST



(b) CelebA



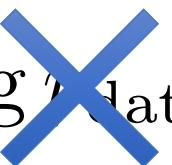
(c) CIFAR-10

Langevin sampling process

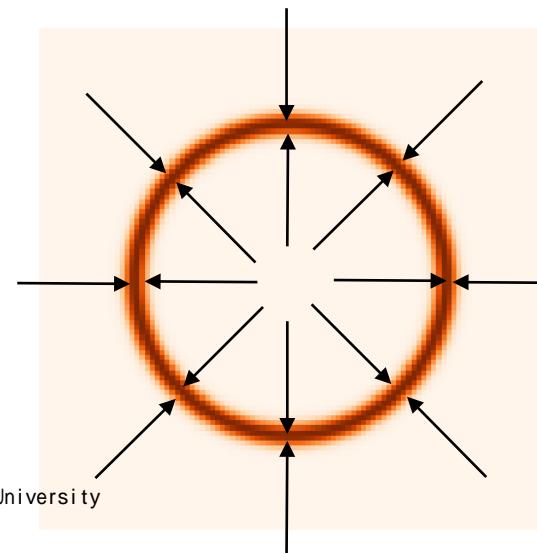
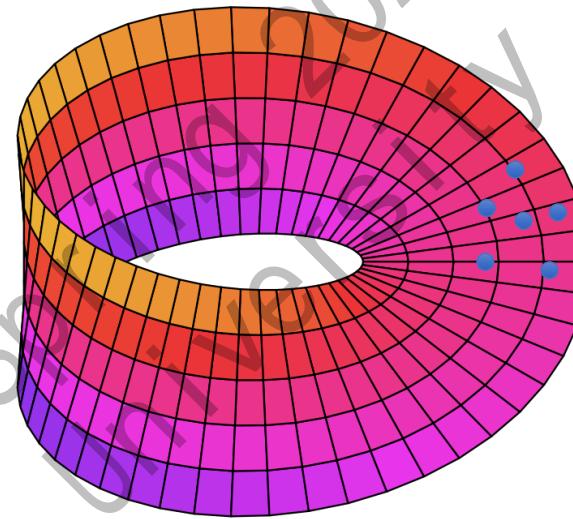
Why does it fail in practice???

# Pitfall 1: manifold hypothesis

- Manifold hypothesis.

$$\nabla_{\mathbf{x}} \log \gamma_{\text{data}}(\mathbf{x})$$


- Data score is undefined.
- Example
  - The data distribution is a ring
  - What is the score function like?
  - **What about the real-world data?**



# Pitfall 1: manifold hypothesis

- A toy example for the manifold hypothesis
  - Fitting the data with a low-dimensional linear manifold (PCA)

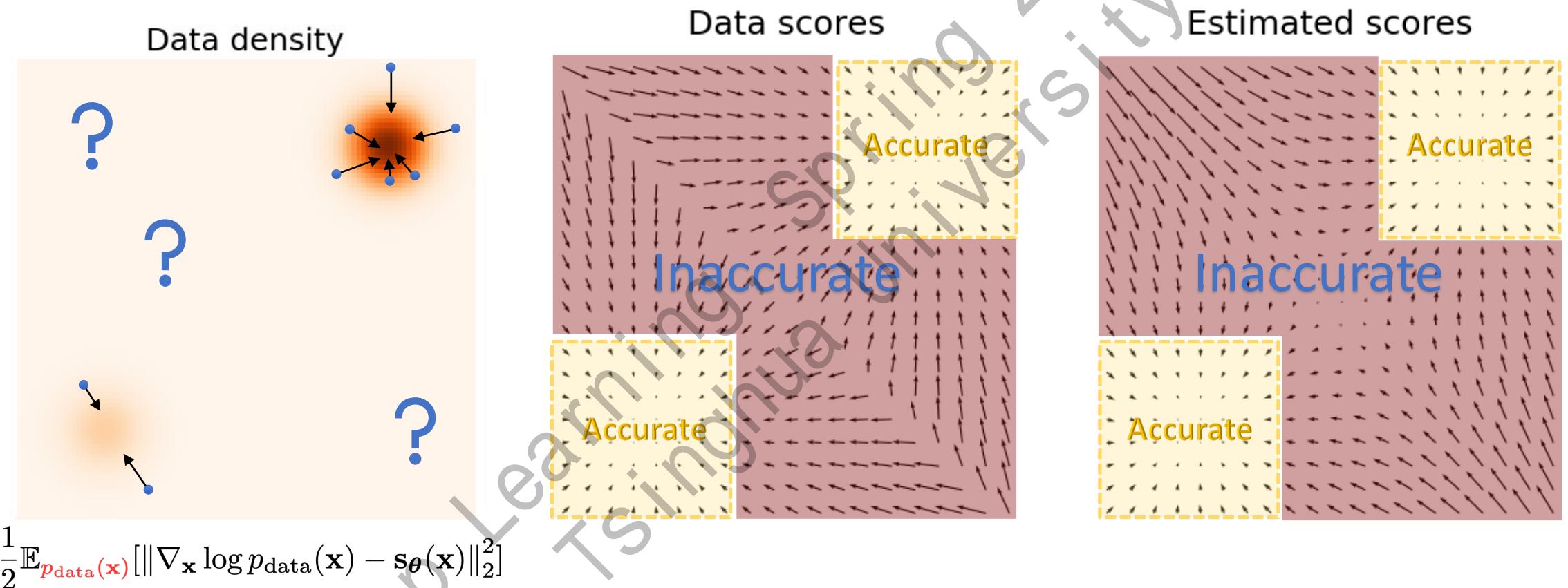


Real-world data have a small intrinsic dimension!

## Pitfall 2: challenge in **low data density regions**

- Let's assume a well defined score function over the entire space

# Pitfall 2: challenge in low data density regions



**Langevin MCMC will have trouble exploring low density regions**

## Pitfall 3: slow mixing of Langevin dynamics **between data modes**

- Let's further assume that we have learned accurate score functions!
- We may still have issue when  $p(x)$  is a multi-modal distribution

## Pitfall 3: slow mixing of Langevin dynamics between data modes

- Suppose the data distribution has two disjoint modes:

$$p_{\text{data}}(\mathbf{x}) = \pi p_1(\mathbf{x}) + (1 - \pi)p_2(\mathbf{x})$$

$$\mathcal{A} \cap \mathcal{B} = \emptyset \quad p_{\text{data}}(\mathbf{x}) = \begin{cases} \pi p_1(\mathbf{x}), & \mathbf{x} \in \mathcal{A} \\ (1 - \pi)p_2(\mathbf{x}), & \mathbf{x} \in \mathcal{B} \end{cases}$$

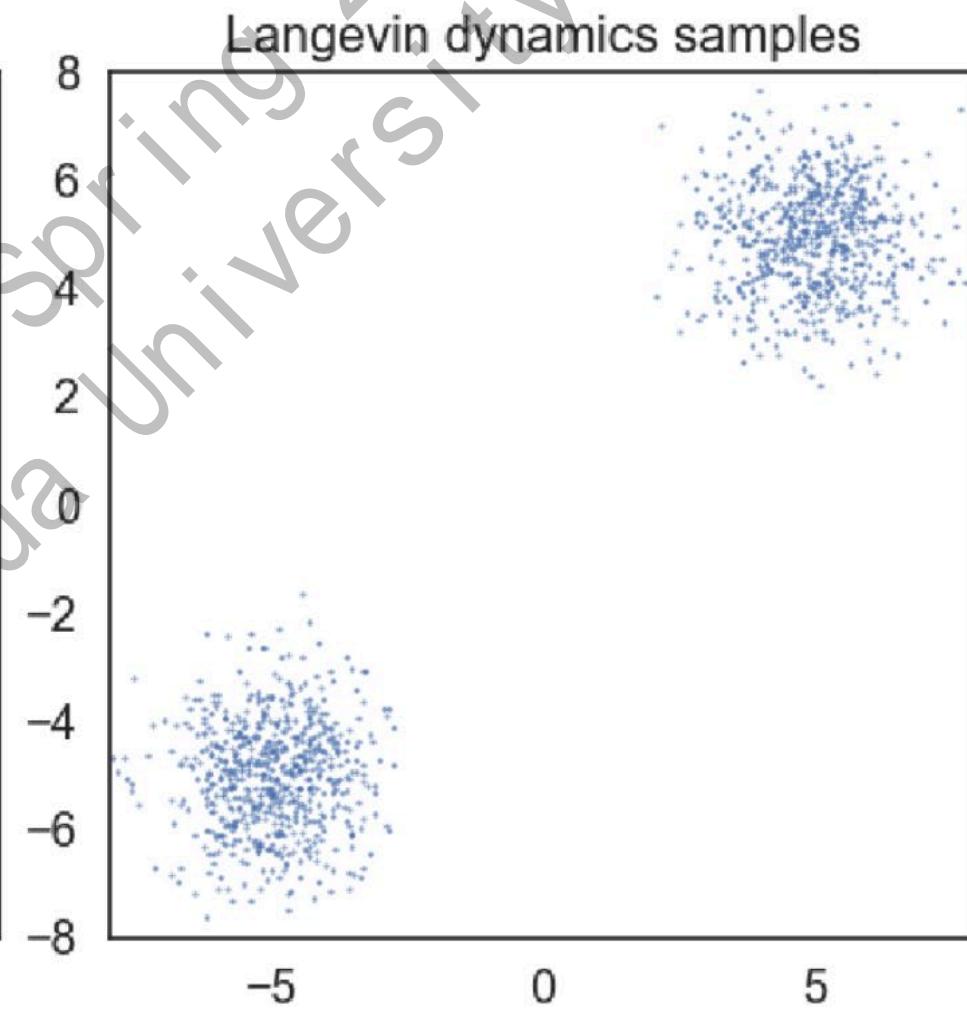
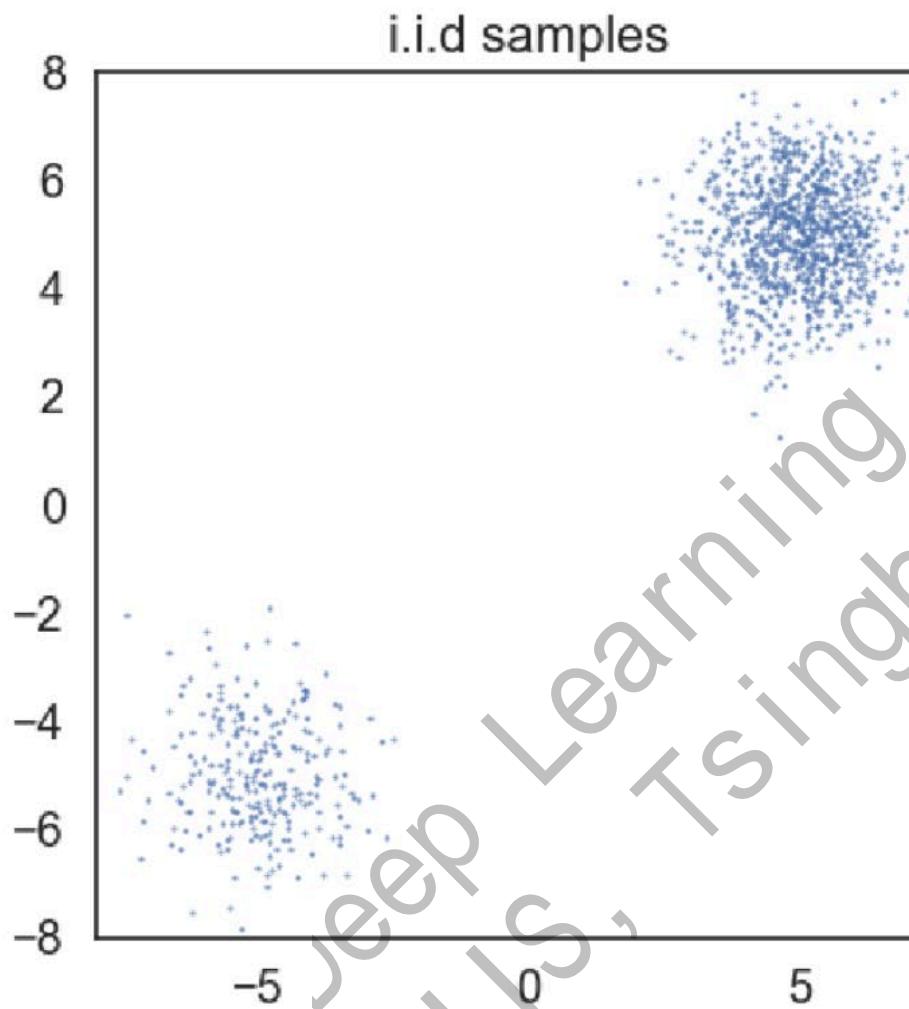
- Data score function:

$$\nabla_{\mathbf{x}} \log p_{\text{data}}(\mathbf{x}) = \begin{cases} \nabla_{\mathbf{x}}[\log \pi + \log p_1(\mathbf{x})], & \mathbf{x} \in \mathcal{A} \\ \nabla_{\mathbf{x}}[\log(1 - \pi) + \log p_2(\mathbf{x})], & \mathbf{x} \in \mathcal{B} \end{cases}$$

$$= \begin{cases} \nabla_{\mathbf{x}} \log p_1(\mathbf{x}), & \mathbf{x} \in \mathcal{A} \\ \nabla_{\mathbf{x}} \log p_2(\mathbf{x}), & \mathbf{x} \in \mathcal{B} \end{cases}$$

- The score function has no dependence on the mode weighting  $\pi$  at all!
- Langevin sampling will not reflect  $\pi$

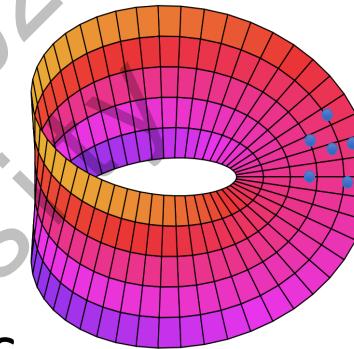
# Pitfall 3: slow mixing of Langevin dynamics between data modes



# Pitfalls

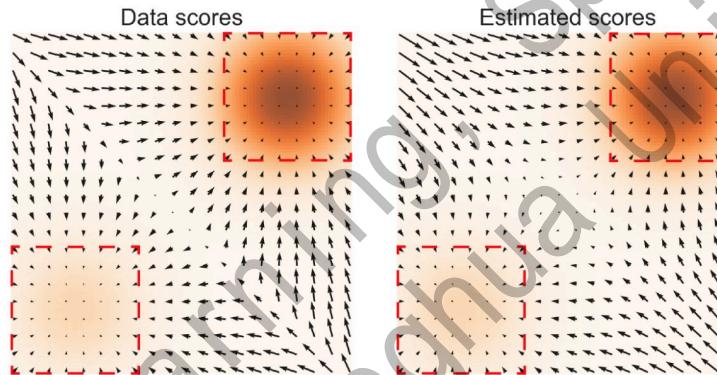
- Manifold hypothesis. Data score is undefined.

$$\nabla_x \log p_{\text{data}}(x)$$

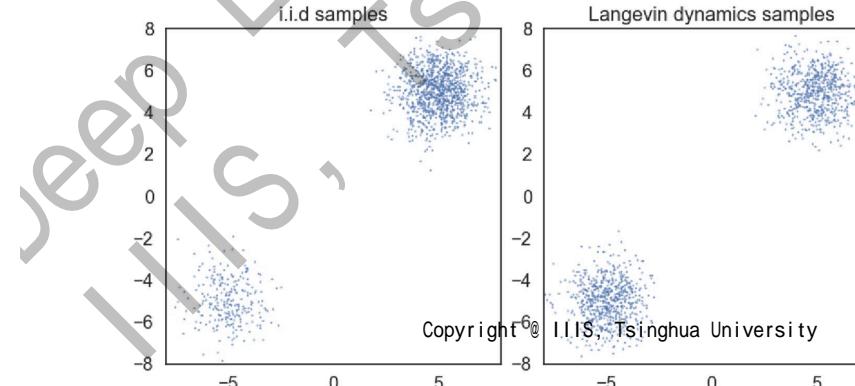


Data points

- Score matching fails in low data density regions

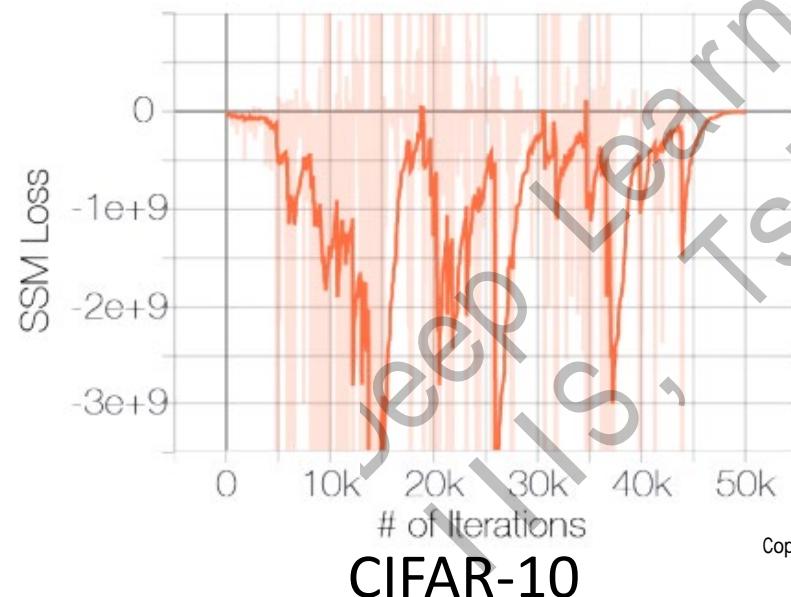
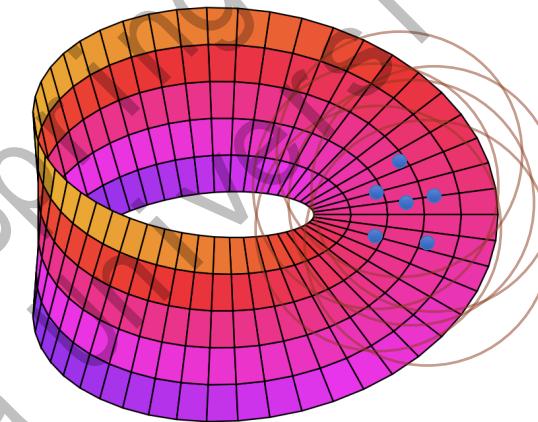


- Langevin dynamics fail to weight different modes correctly



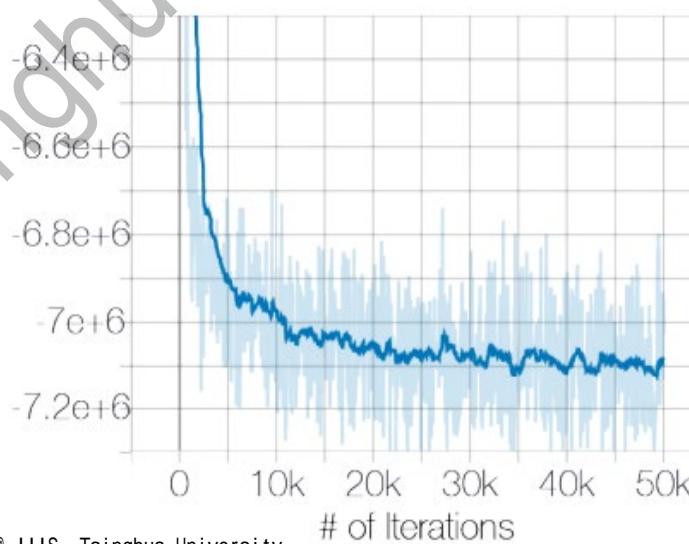
# Gaussian perturbation

- The solution to all pitfalls: **Gaussian perturbation!**
- Manifold + noise
- Score matching on noisy data.

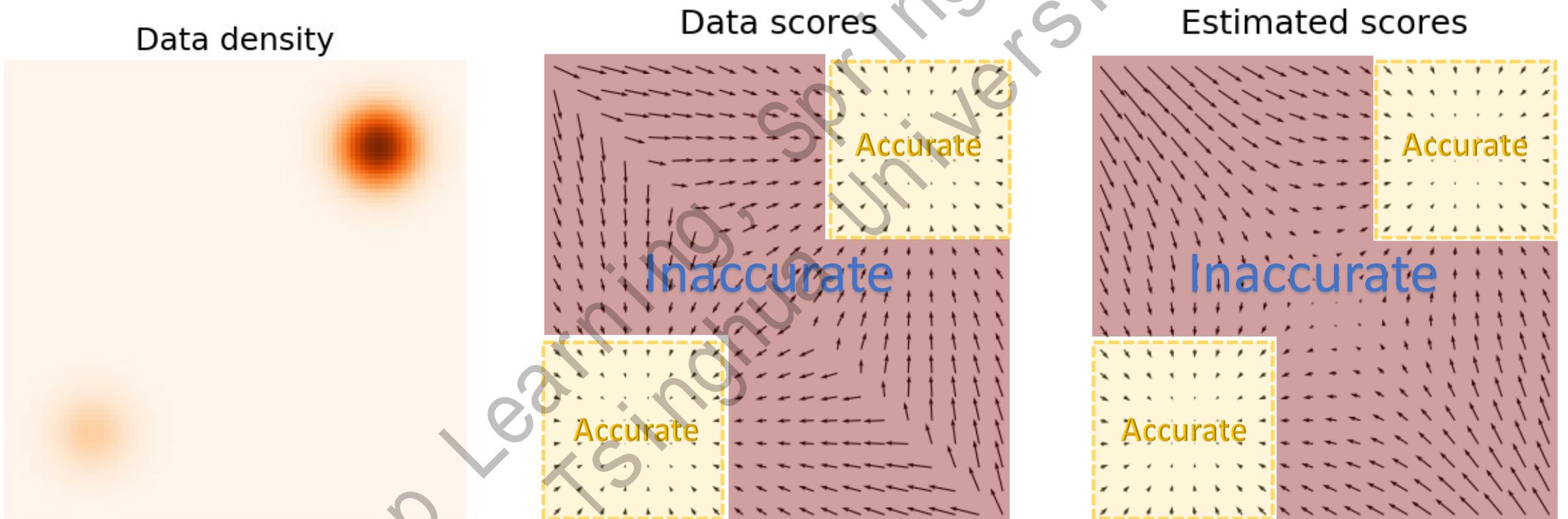


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CIFAR-10

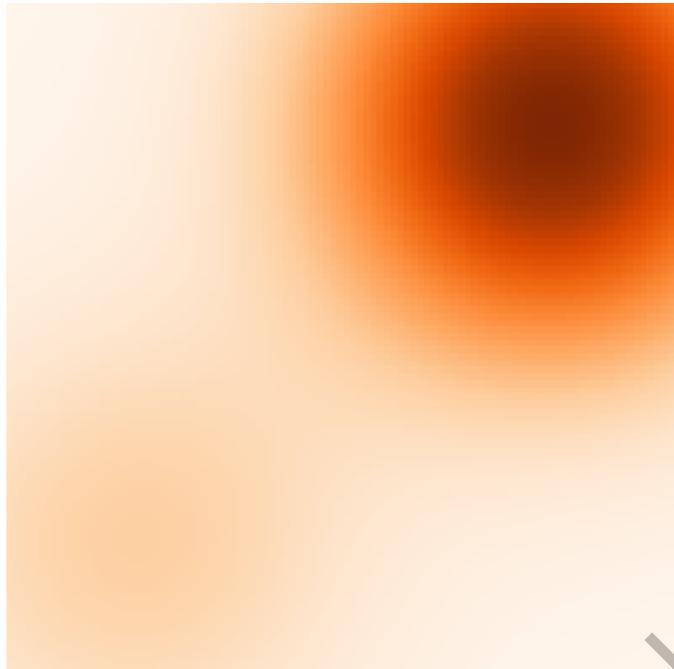
 $\mathcal{N}(0; 0.0001)$

# Challenge in low data density regions

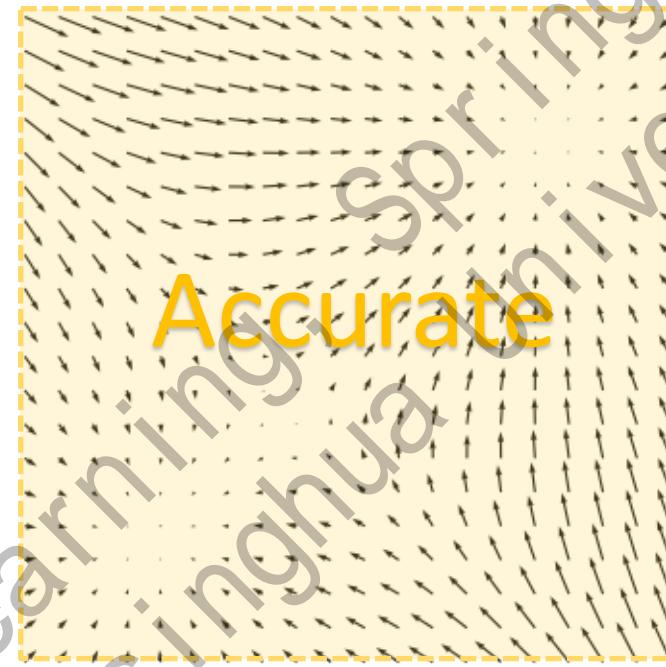


# Adding noise to data

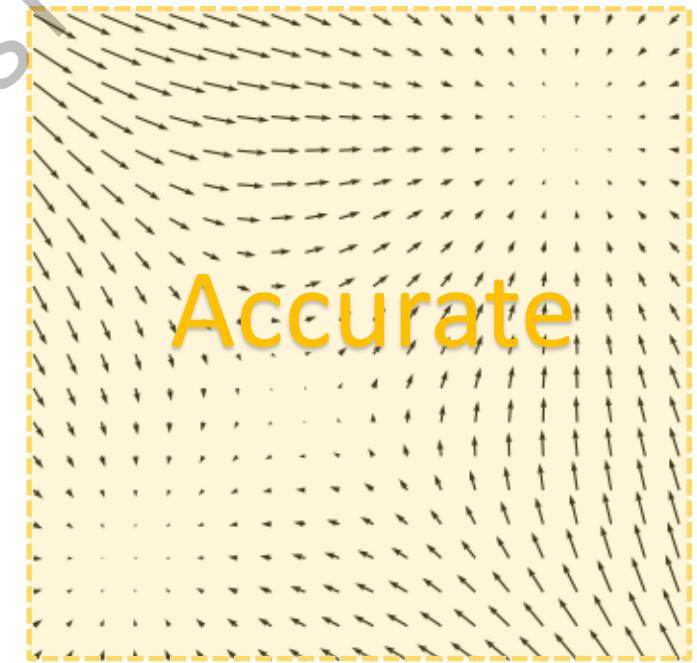
Perturbed density



Perturbed scores



Estimated scores



Provide useful directional information  
for Langevin MCMC.

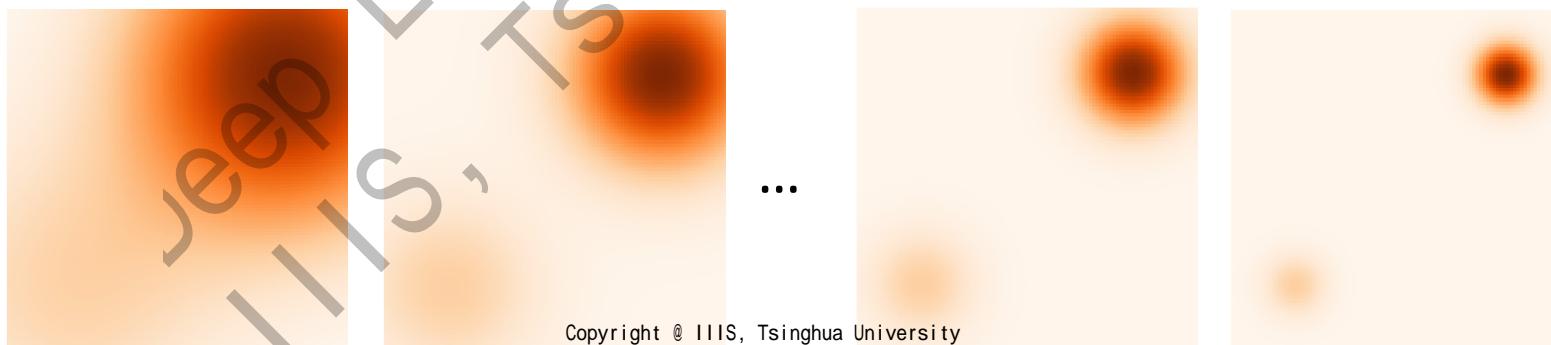
# Multi-scale Noise Perturbation

- Trade-off

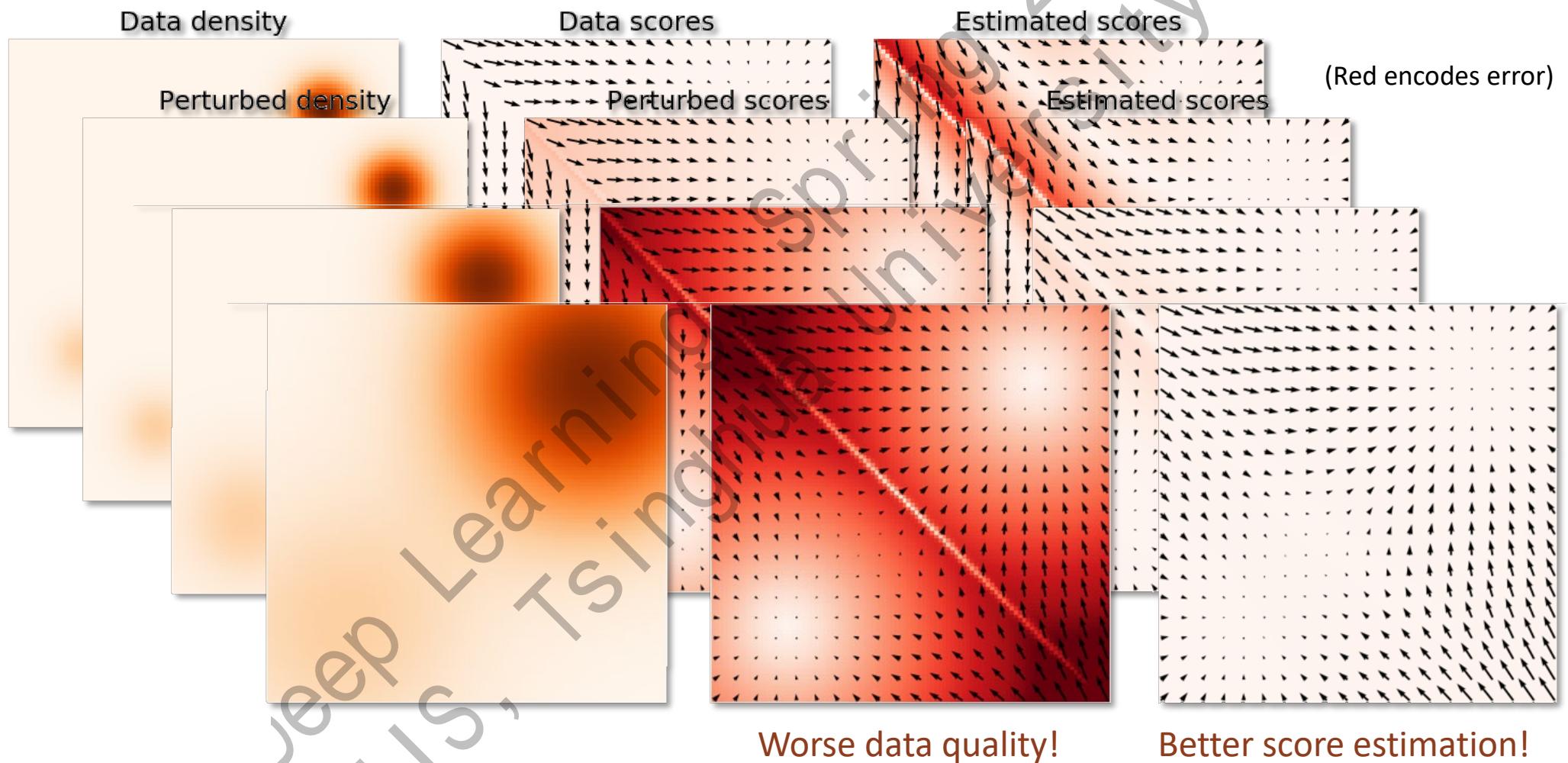


- Multi-scale noise perturbations.

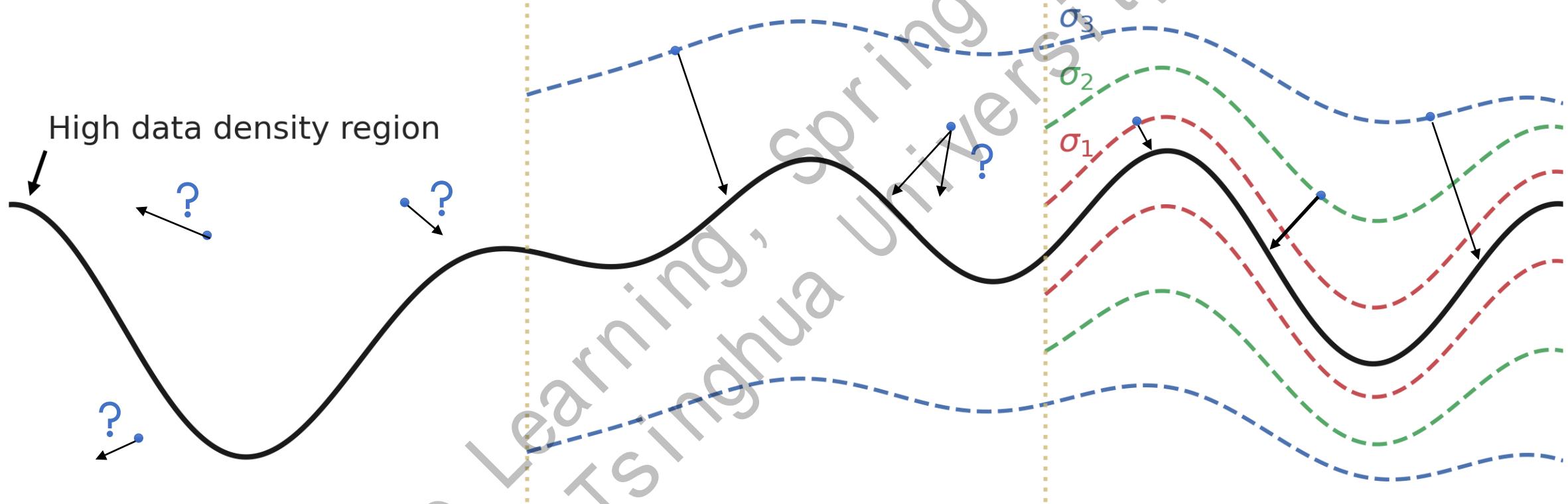
$$\sigma_1 > \sigma_2 > \dots > \sigma_{L-1} > \sigma_L$$



# Trading off Data Quality and Estimation Accuracy

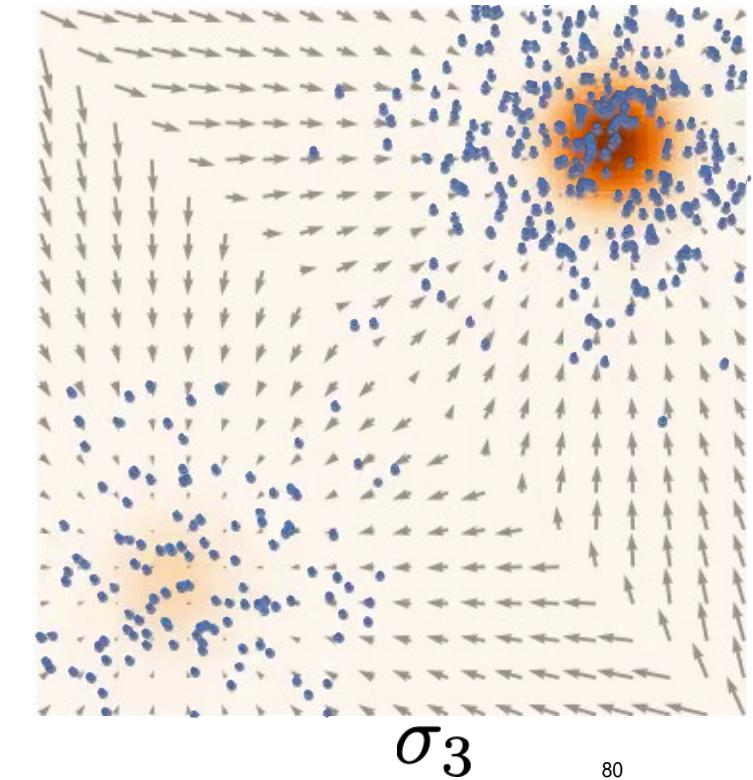
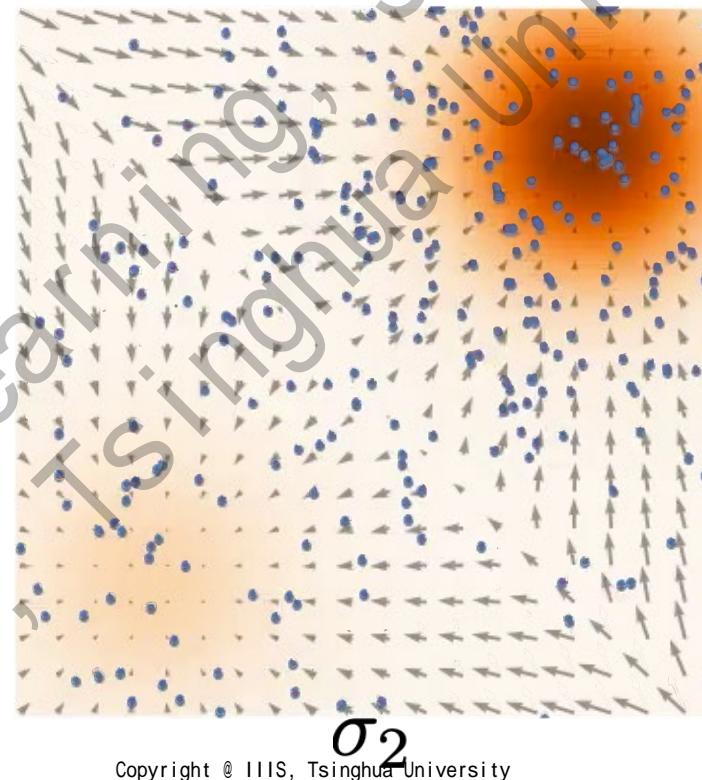
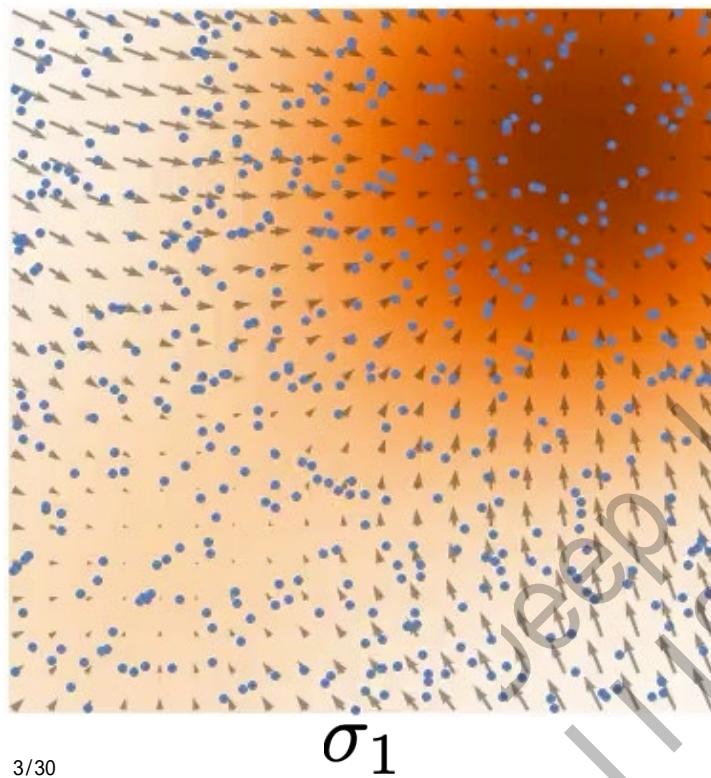


# Using multiple noise scales



# Annealed Langevin Dynamics: Joint Scores to Samples

- Sample using  $\sigma_1, \sigma_2, \dots, \sigma_L$  sequentially with Langevin dynamics.
- Anneal down the noise level.
- Samples used as initialization for the next level.



# Annealed Langevin dynamics

---

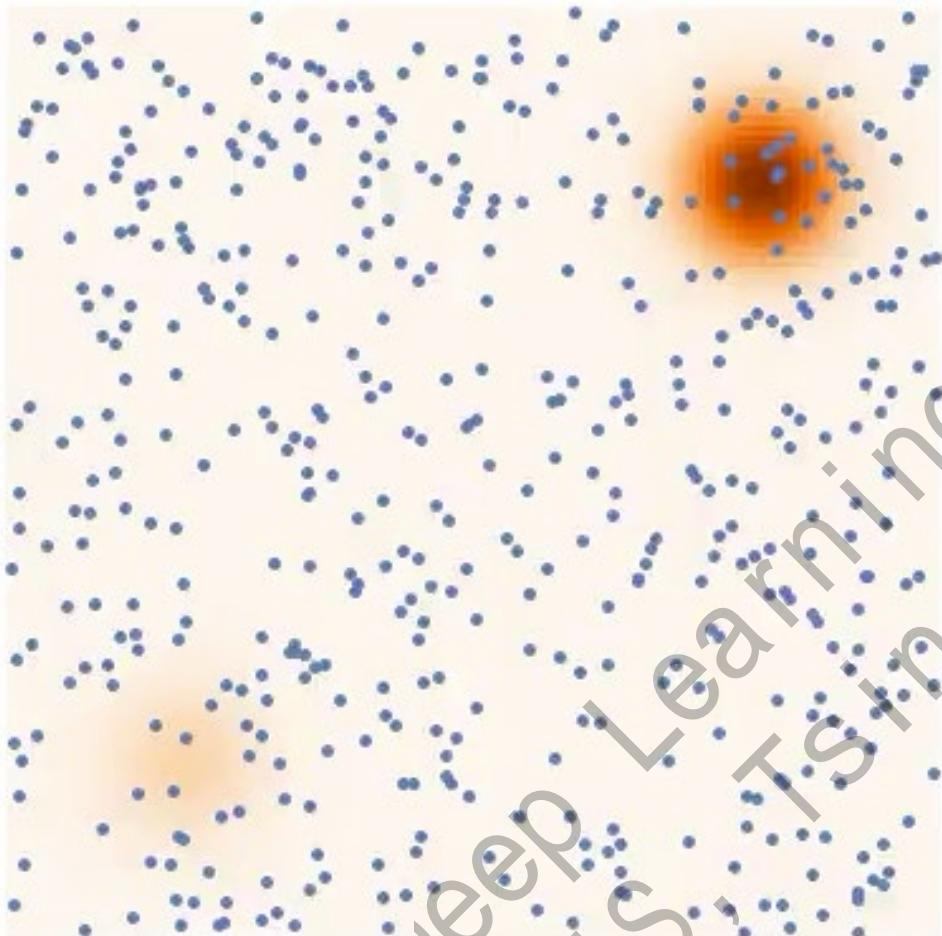
## Algorithm 1 Annealed Langevin dynamics.

---

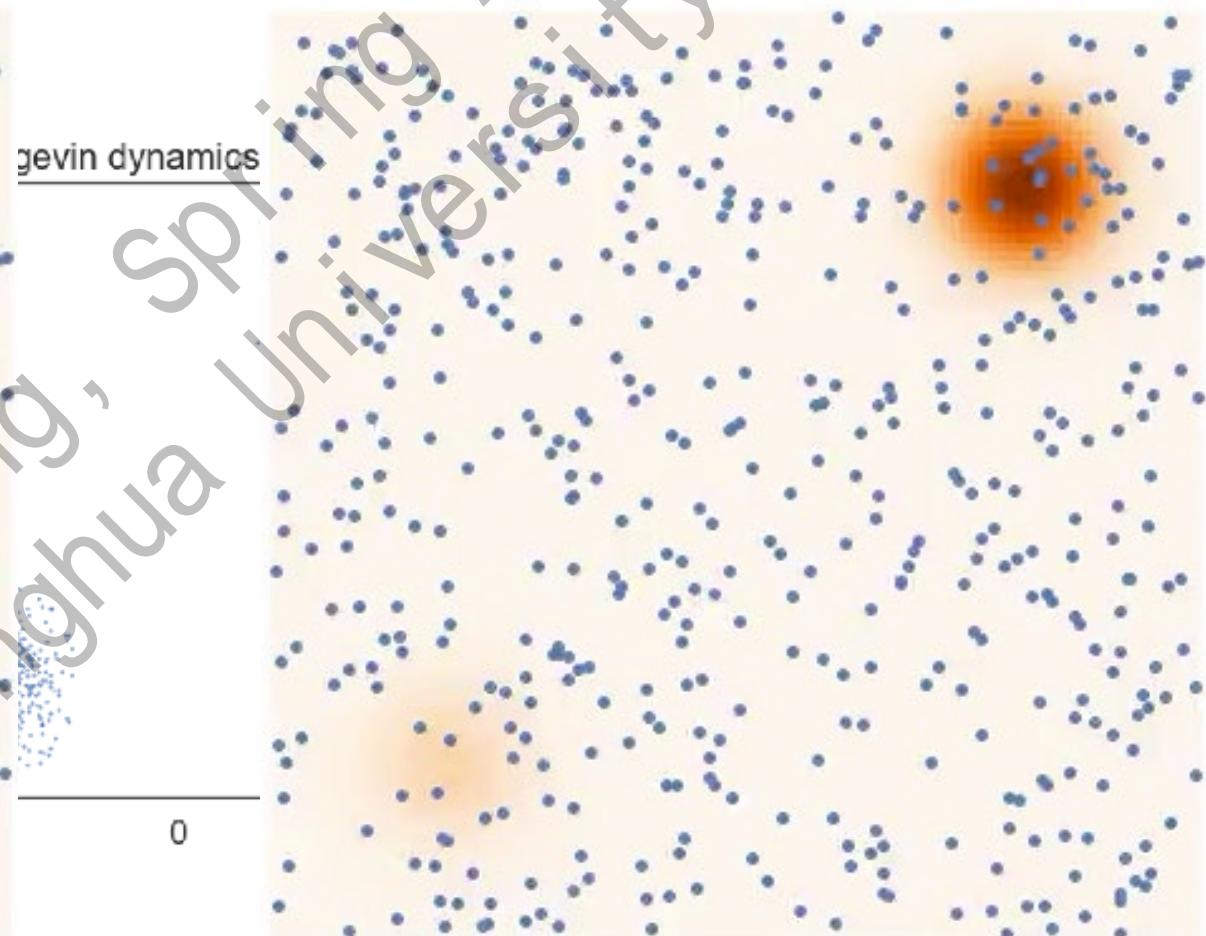
**Require:**  $\{\sigma_i\}_{i=1}^L, \epsilon, T$ .

```
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  $\tilde{\mathbf{x}}_T$ 
```

# Comparison to the vanilla Langevin dynamics

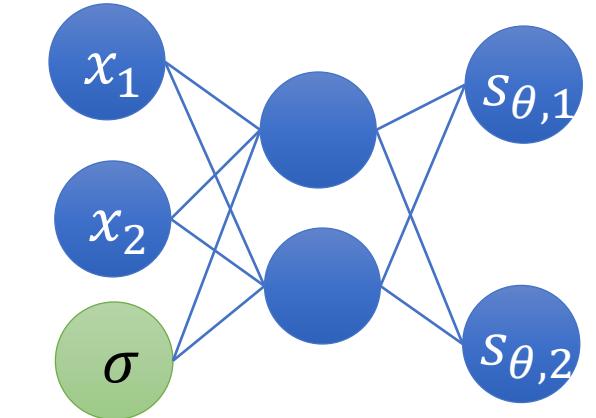
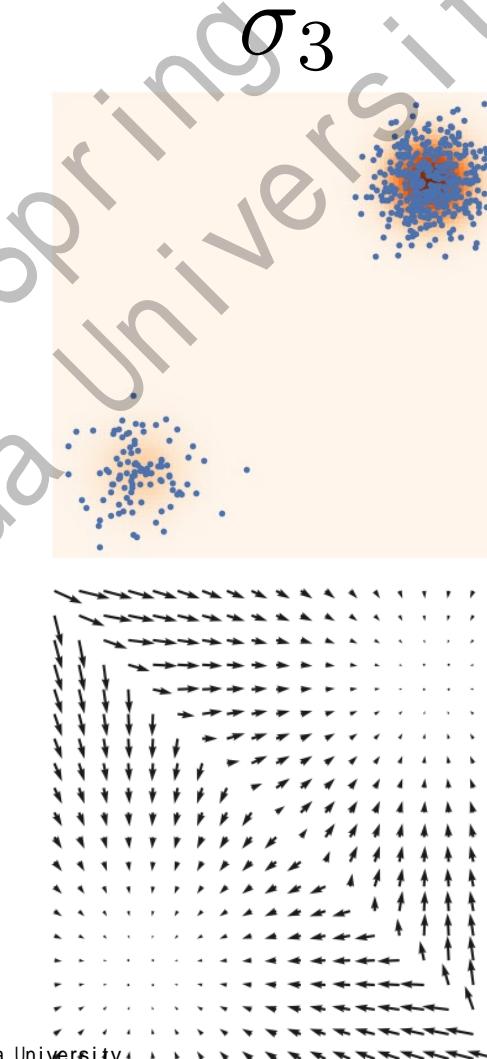
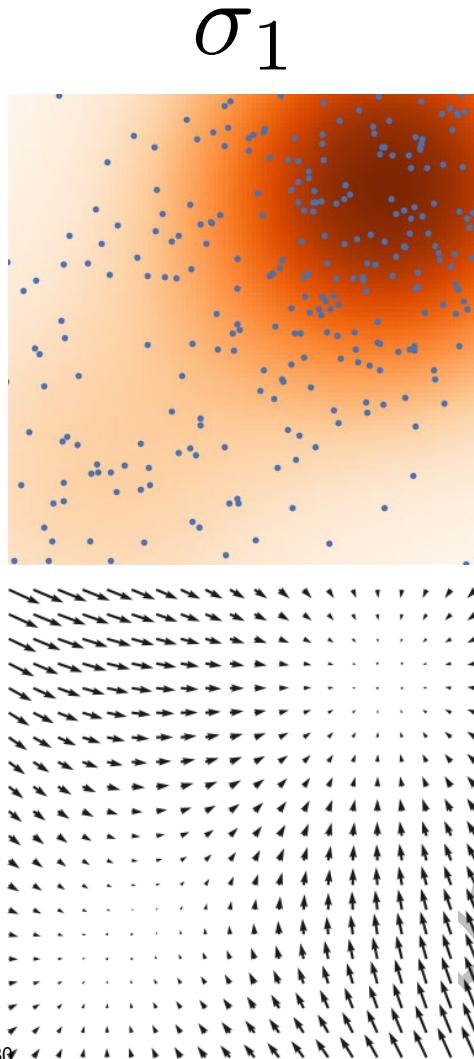


Langevin dynamics



Annealed Langevin dynamics

# Joint Score Estimation via Noise Conditional Score Networks



Noise Conditional  
Score Network  
(NCSN)

# Training noise conditional score networks

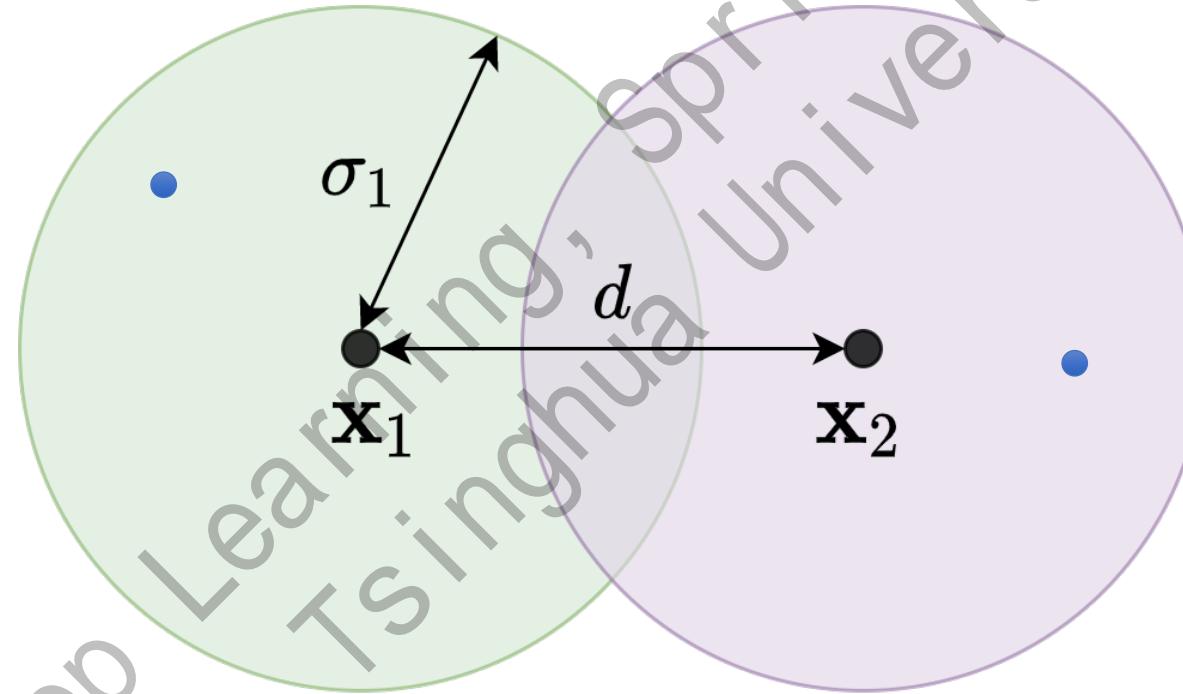
- Weighted combination of denoising score matching losses
  - Given the noise levels  $\sigma_1 \dots \sigma_L$

$$\begin{aligned}
 & \frac{1}{L} \sum_{i=1}^L \lambda(\sigma_i) E_{q_{\sigma_i}(\mathbf{x})} [\|\nabla_{\mathbf{x}} \log q_{\sigma_i}(\mathbf{x}) - \mathbf{s}_{\theta}(\mathbf{x}, \sigma_i)\|_2^2] \\
 &= \frac{1}{L} \sum_{i=1}^L \lambda(\sigma_i) E_{\mathbf{x} \sim p_{\text{data}}, \mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} [\|\nabla_{\tilde{\mathbf{x}}} \log q_{\sigma_i}(\tilde{\mathbf{x}} \mid \mathbf{x}) - \mathbf{s}_{\theta}(\tilde{\mathbf{x}}, \sigma_i)\|_2^2] + \text{const.} \\
 &= \frac{1}{L} \sum_{i=1}^L \lambda(\sigma_i) E_{\mathbf{x} \sim p_{\text{data}}, \mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[ \left\| \mathbf{s}_{\theta}(\mathbf{x} + \sigma_i \mathbf{z}, \sigma_i) + \frac{\mathbf{z}}{\sigma_i} \right\|_2^2 \right] + \text{const.}
 \end{aligned}$$

# Choosing noise scales

- Maximum noise scale

$\sigma_1 \approx$  maximum pairwise distance between datapoints



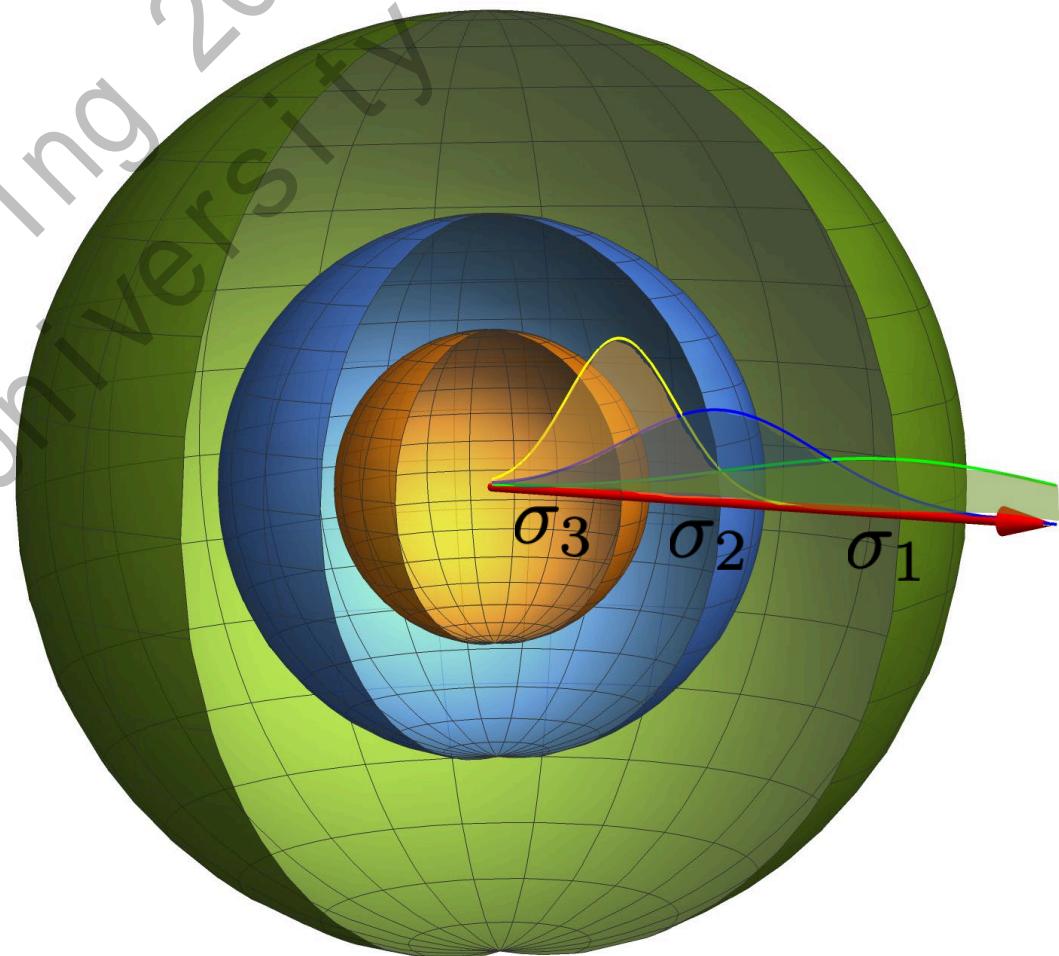
- Minimum noise scale:  $\sigma_L$  should be sufficiently small to control the noise in final samples.

# Choosing noise scales

- **Key intuition:** adjacent noise scales should have sufficient overlap to facilitate transitioning across noise scales in annealed Langevin dynamics.
- A geometric progression with sufficient length.

$$\sigma_1 > \sigma_2 > \sigma_3 > \dots > \sigma_{L-1} > \sigma_L$$

$$\frac{\sigma_1}{\sigma_2} = \frac{\sigma_2}{\sigma_3} = \dots = \frac{\sigma_{L-1}}{\sigma_L}$$



# Choosing the weighting function

- Weighted combination of denoising score matching losses

$$\frac{1}{L} \sum_{i=1}^L \lambda(\sigma_i) E_{\mathbf{x} \sim p_{\text{data}}, \mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[ \left\| \mathbf{s}_\theta(\mathbf{x} + \sigma_i \mathbf{z}, \sigma_i) + \frac{\mathbf{z}}{\sigma_i} \right\|_2^2 \right]$$

- How to choose the weighting function  $\lambda : \mathbb{R}_{>0} \rightarrow \mathbb{R}_{>0}$
- Goal:** balancing different score matching losses in the sum  $\rightarrow \lambda(\sigma_i) = \sigma_i^2$

$$\begin{aligned} & \frac{1}{L} \sum_{i=1}^L \sigma_i^2 E_{\mathbf{x} \sim p_{\text{data}}, \mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[ \left\| \mathbf{s}_\theta(\mathbf{x} + \sigma_i \mathbf{z}, \sigma_i) + \frac{\mathbf{z}}{\sigma_i} \right\|_2^2 \right] \\ &= \frac{1}{L} \sum_{i=1}^L E_{\mathbf{x} \sim p_{\text{data}}, \mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[ \left\| \sigma_i \mathbf{s}_\theta(\mathbf{x} + \sigma_i \mathbf{z}, \sigma_i) + \mathbf{z} \right\|_2^2 \right] \\ &= \frac{1}{L} \sum_{i=1}^L E_{\mathbf{x} \sim p_{\text{data}}, \mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[ \left\| \boldsymbol{\epsilon}_\theta(\mathbf{x} + \sigma_i \mathbf{z}, \sigma_i) + \mathbf{z} \right\|_2^2 \right] \quad [ \boldsymbol{\epsilon}_\theta(\cdot, \sigma_i) := \sigma_i \mathbf{s}_\theta(\cdot, \sigma_i) ] \end{aligned}$$

# Training noise conditional score networks

- Sample a mini-batch of datapoints  $\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\} \sim p_{\text{data}}$
- Sample a mini-batch of noise scale indices

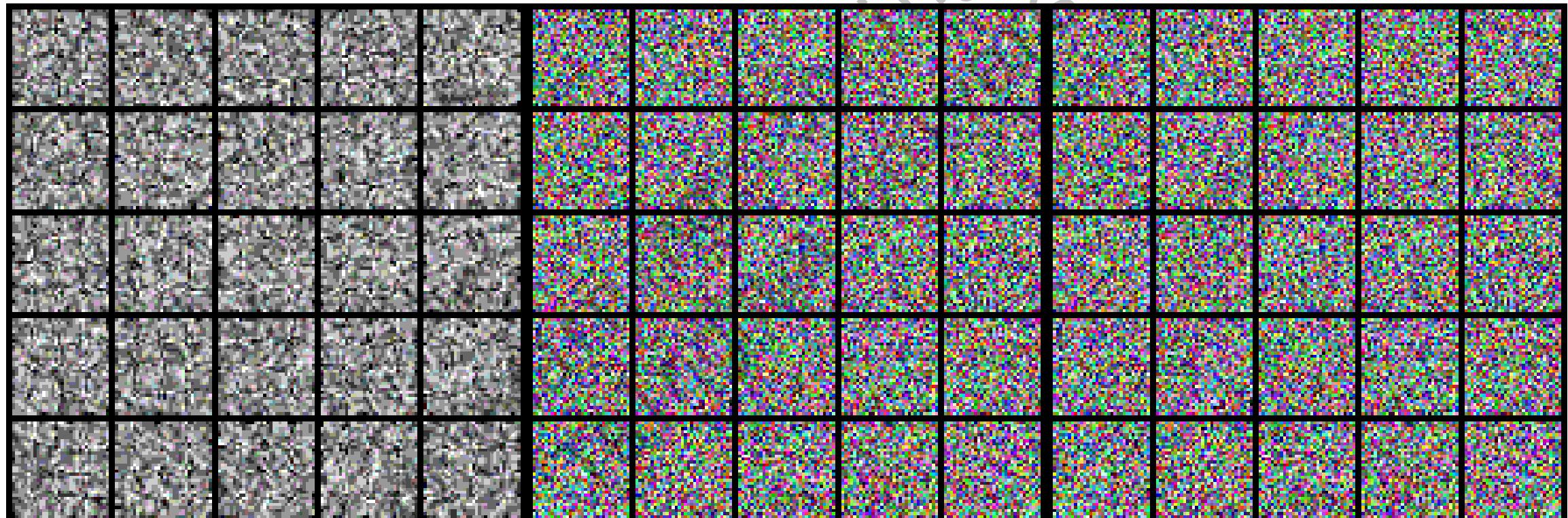
$$\{i_1, i_2, \dots, i_n\} \sim \mathcal{U}\{1, 2, \dots, L\}$$

- Sample a mini-batch of Gaussian noise  $\{\mathbf{z}_1, \mathbf{z}_2, \dots, \mathbf{z}_n\} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$
- Estimate the weighted mixture of score matching losses

$$\frac{1}{n} \sum_{k=1}^n \left[ \|\sigma_{i_k} s_{\theta}(\mathbf{x}_k + \sigma_{i_k} \mathbf{z}_k, \sigma_{i_k}) + \mathbf{z}_k\|_2^2 \right]$$

- Stochastic gradient descent
- As efficient as training one single non-conditional score-based model

# Experiments: Sampling



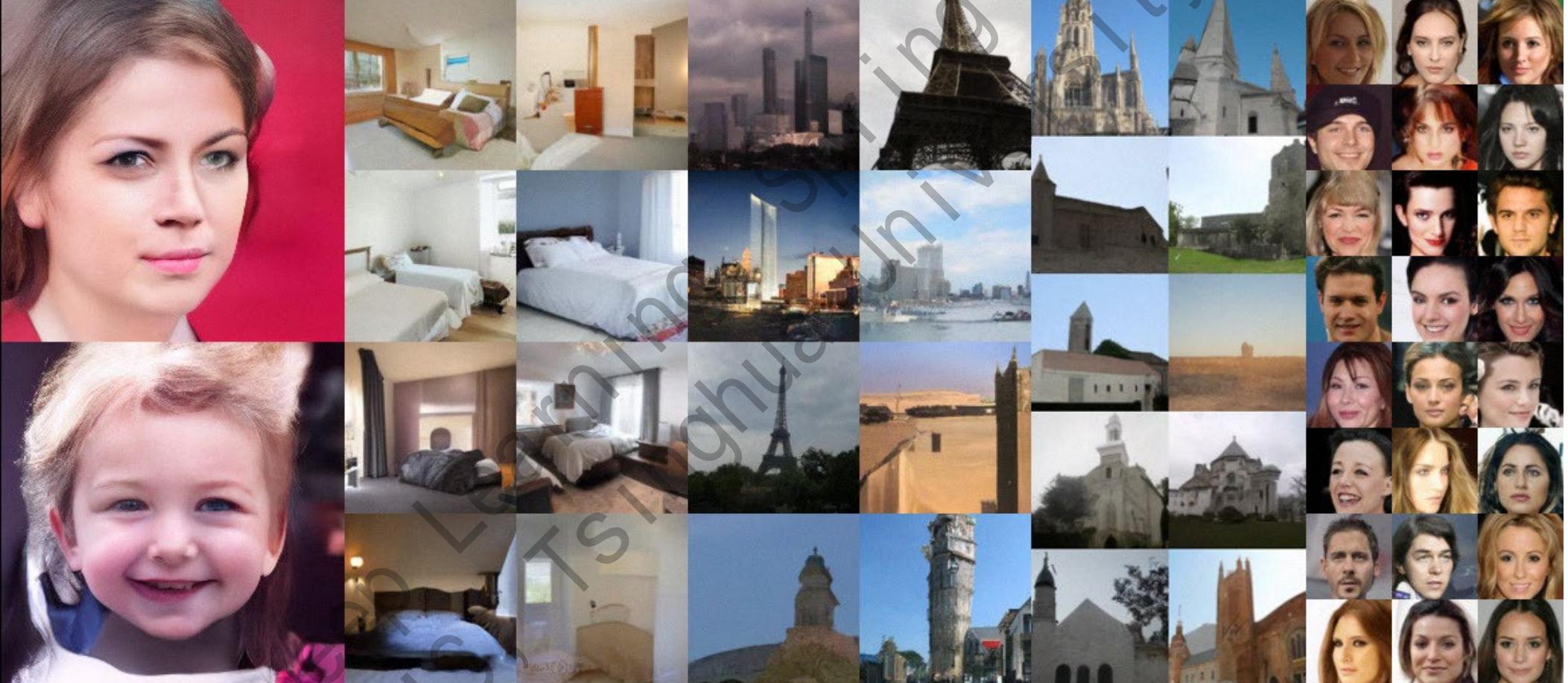
# Experiments: Sampling

---

Model	Inception	FID
<b>CIFAR-10 Unconditional</b>		
PixelCNN [59]	4.60	65.93
PixelIQN [42]	5.29	49.46
EBM [12]	6.02	40.58
WGAN-GP [18]	$7.86 \pm .07$	36.4
MoLM [45]	$7.90 \pm .10$	<b>18.9</b>
SNGAN [36]	$8.22 \pm .05$	21.7
ProgressiveGAN [25]	$8.80 \pm .05$	-
<b>NCSN (Ours)</b>	<b><math>8.87 \pm .12</math></b>	25.32

---

# High Resolution Image Generation



# Comparison between NCSN and DDPM

- NCSN

- Learning: 
$$\frac{1}{L} \sum_{i=1}^L E_{\mathbf{x} \sim p_{\text{data}}, \mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[ \|\epsilon_{\theta}(\mathbf{x} + \sigma_i \mathbf{z}, \sigma_i) + \mathbf{z}\|_2^2 \right] \quad [\epsilon_{\theta}(\cdot, \sigma_i) := \sigma_i s_{\theta}(\cdot, \sigma_i)]$$
- Inference: 
$$\tilde{\mathbf{x}}_t \leftarrow \tilde{\mathbf{x}}_{t-1} + \frac{\alpha_i}{2} \mathbf{s}_{\theta}(\tilde{\mathbf{x}}_{t-1}, \sigma_i) + \sqrt{\alpha_i} \mathbf{z}_t$$

- DDPM

- NCSN with a few enhancements for training stabilities
  - More discussions can be found in the original paper (your homework ☺)

---

**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
      
$$\nabla_{\theta} \|\epsilon - \epsilon_{\theta}(\sqrt{\alpha_t} \mathbf{x}_0 + \sqrt{1 - \alpha_t} \epsilon, t)\|^2$$

6: until converged
  
```

---

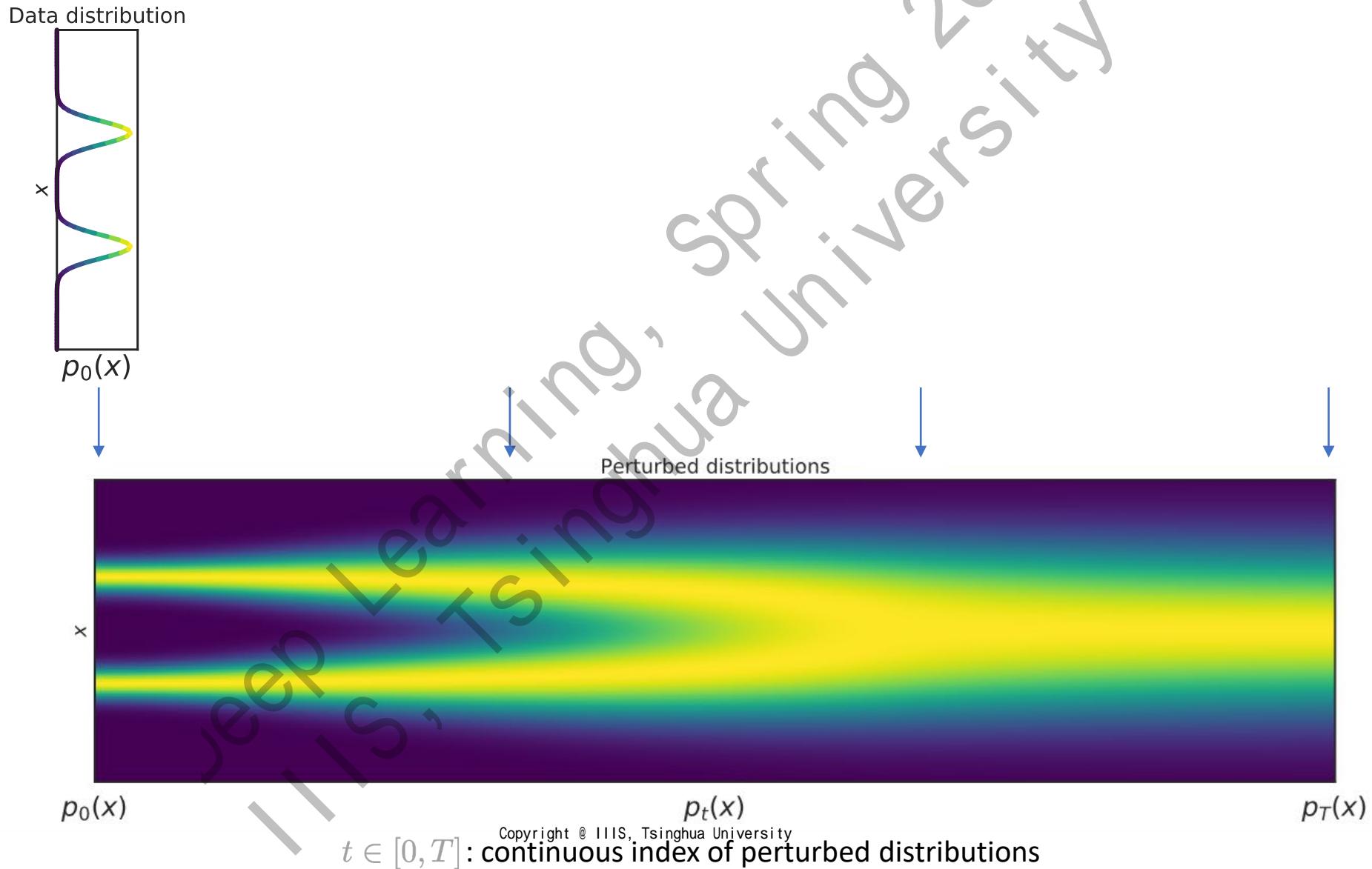
**Algorithm 2** Sampling
 

---

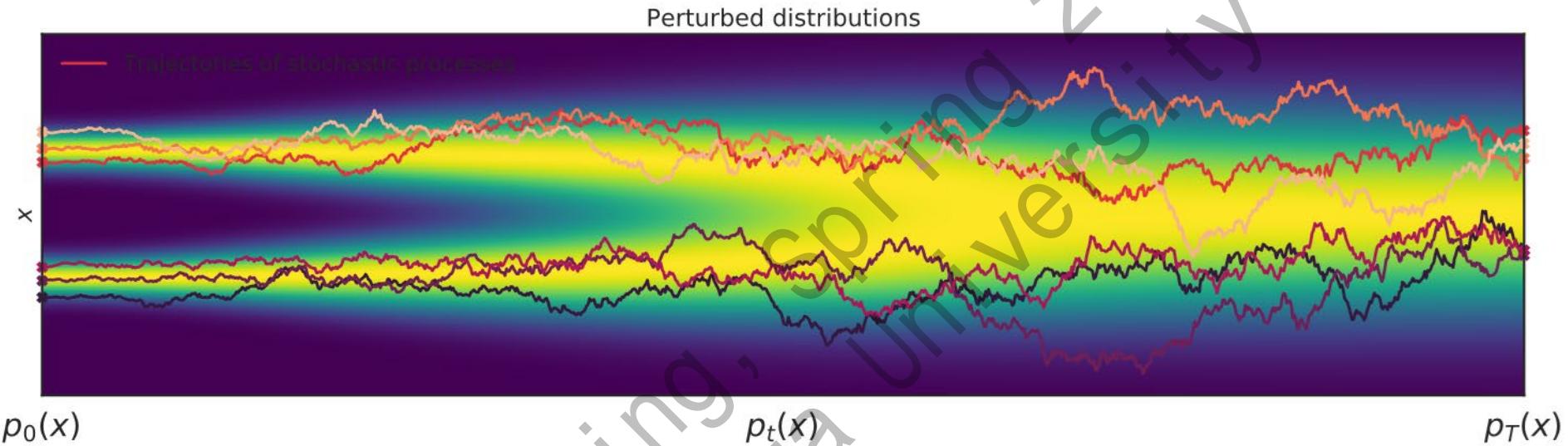
```

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:    $\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left( \mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1 - \alpha_t}} \epsilon_{\theta}(\mathbf{x}_t, t) \right) + \sigma_t \mathbf{z}$ 
5: end for
6: return  $\mathbf{x}_0$ 
  
```

# Using an infinite number of noise scales



# Compact representation of infinite distributions

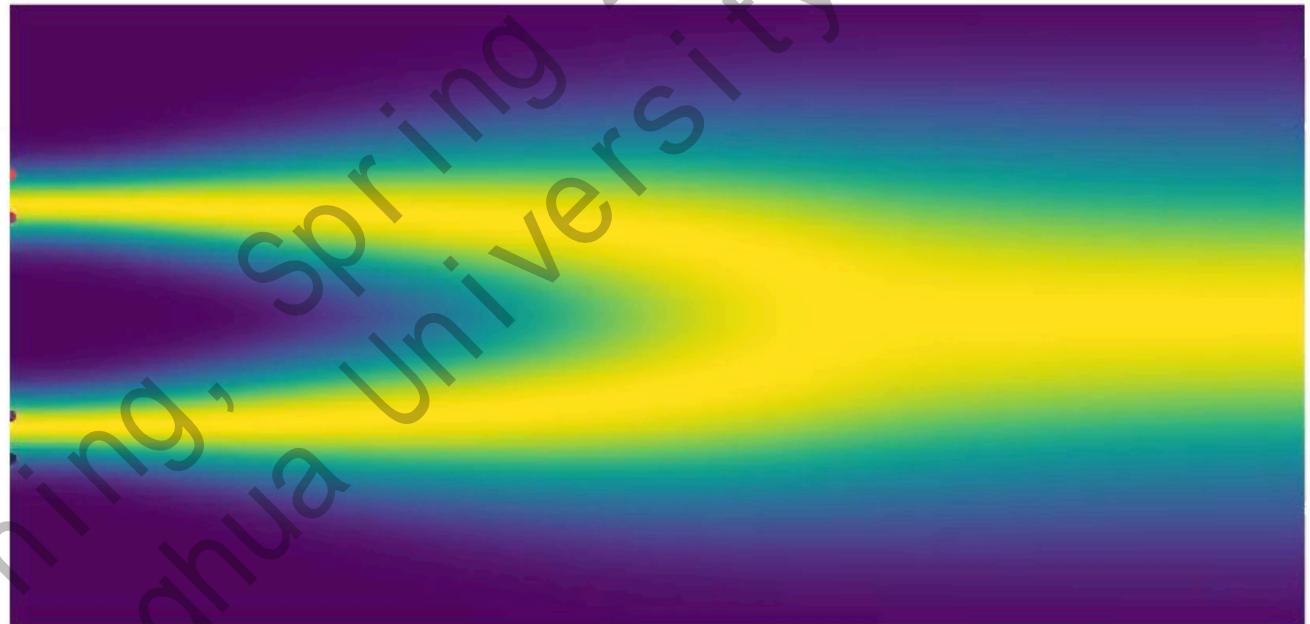


- Stochastic process  $\{\mathbf{x}(t)\}_{t=0}^T \rightarrow$  Marginal probability densities  $\{p_t(\mathbf{x})\}_{t=0}^T$
- Stochastic differential equation:  $d\mathbf{x} = \boxed{\mathbf{f}(\mathbf{x}, t)dt} + \sigma(t) \boxed{d\mathbf{w}}$ 

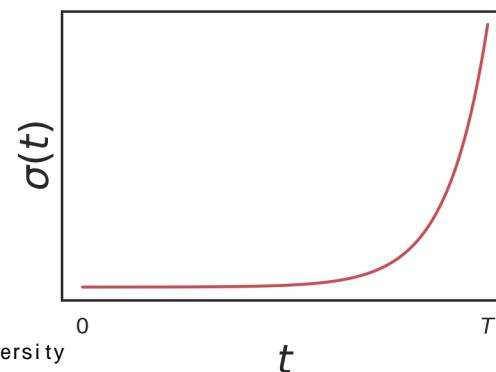
**Deterministic drift** **Infinitesimal white noise**<sup>94</sup>

# Score-based generative modeling via SDEs

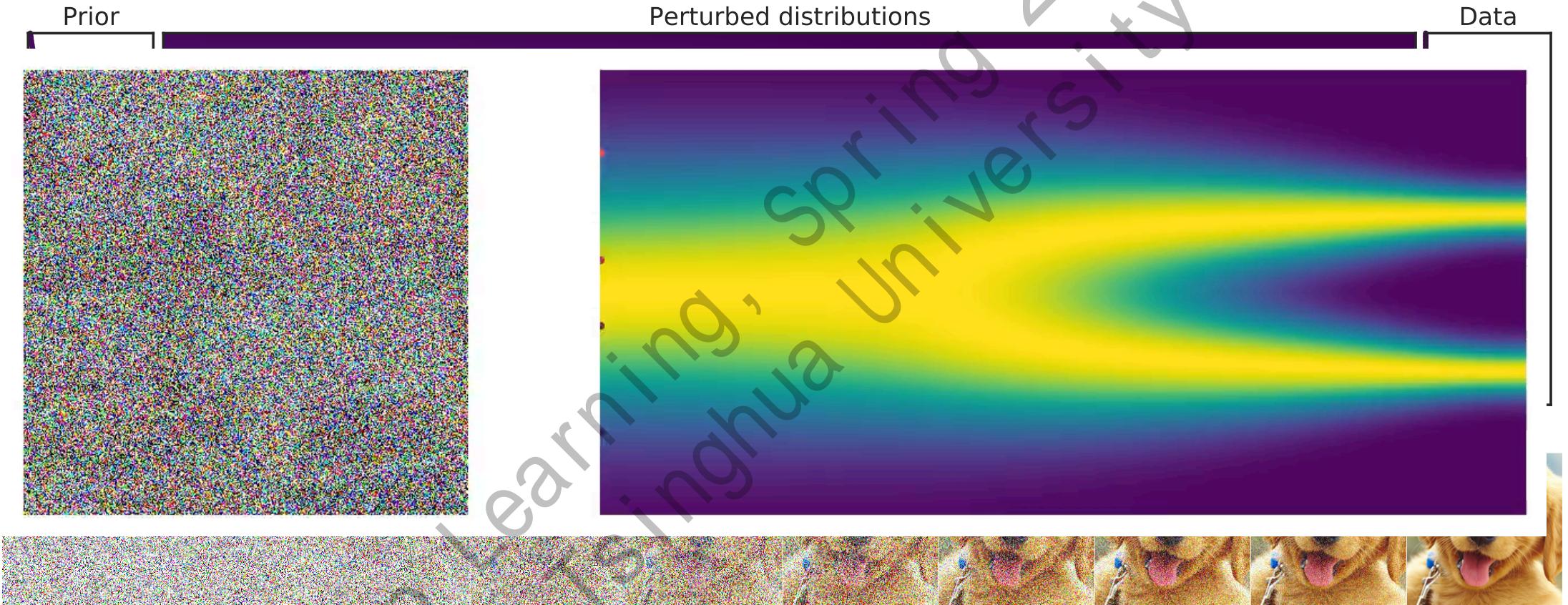
Data



$$d\mathbf{x} = \sigma(t) d\mathbf{w}$$



# Score-based generative modeling via SDEs



$$d\mathbf{x} = \sigma(t)d\mathbf{w} \xrightarrow{\text{Time reversal}} d\mathbf{x} = -\sigma^2(t) \boxed{\nabla_{\mathbf{x}} \log p_t(\mathbf{x})} dt + \sigma(t)d\bar{\mathbf{w}}$$

Score function!

Copyright © IIIS, Tsinghua University

# Score-based generative modeling via SDEs

- Time-dependent score-based model

$$\mathbf{s}_\theta(\mathbf{x}, t) \approx \nabla_{\mathbf{x}} \log p_t(\mathbf{x})$$

- Training:

$$\mathbb{E}_{t \in \mathcal{U}(0, T)} [\lambda(t) \mathbb{E}_{p_t(\mathbf{x})} [\|\nabla_{\mathbf{x}} \log p_t(\mathbf{x}) - \mathbf{s}_\theta(\mathbf{x}, t)\|_2^2]]$$

- Reverse-time SDE

$$d\mathbf{x} = -\sigma^2(t) \mathbf{s}_\theta(\mathbf{x}, t) dt + \sigma(t) d\bar{\mathbf{w}}$$

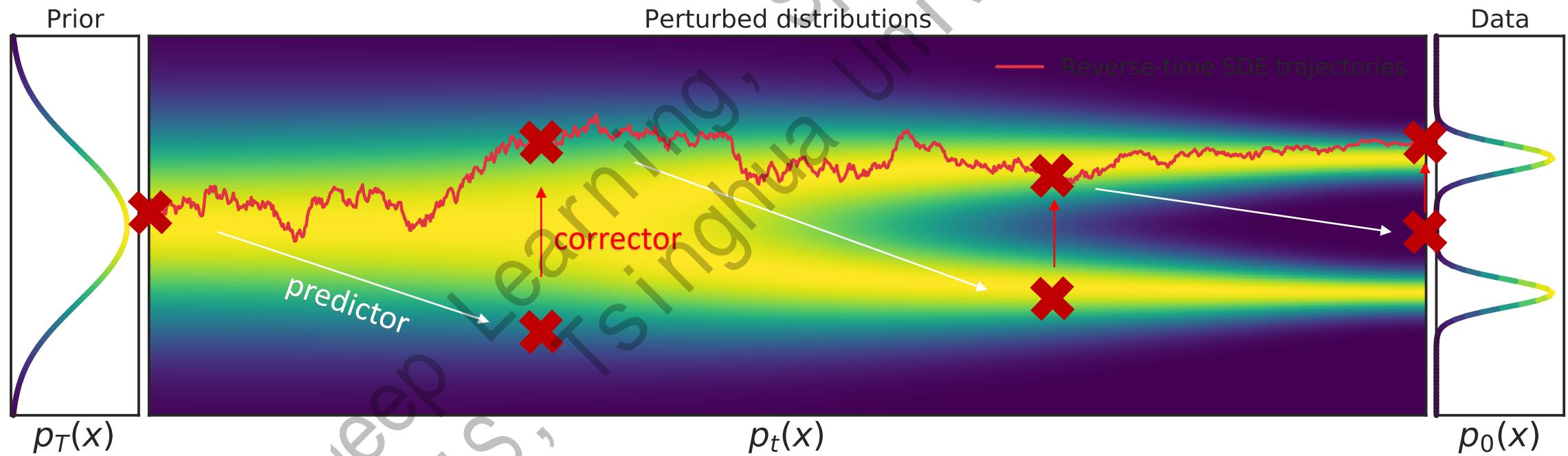
- Sampling: Euler-Maruyama

$$\mathbf{x} \leftarrow \mathbf{x} - \sigma(t)^2 \mathbf{s}_\theta(\mathbf{x}, t) \Delta t + \sigma(t) \mathbf{z} \quad (\mathbf{z} \sim \mathcal{N}(\mathbf{0}, |\Delta t| \mathbf{I}))$$

$$t \leftarrow t + \Delta t$$

# Predictor-Corrector sampling methods

- Predictor-Corrector sampling.
  - **Predictor:** Numerical SDE solver
  - **Corrector:** Score-based MCMC



# Results on predictor-corrector sampling

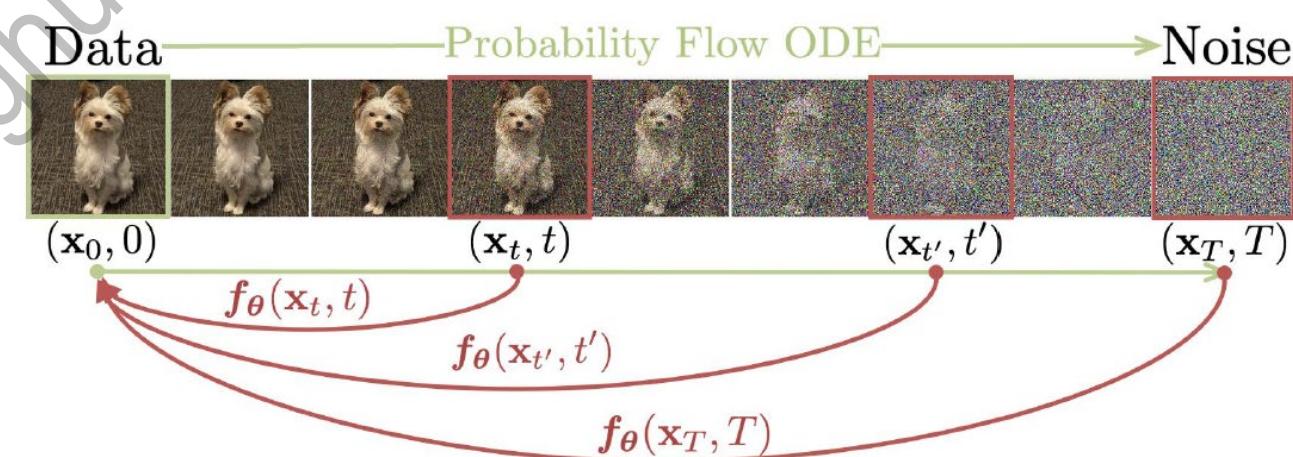
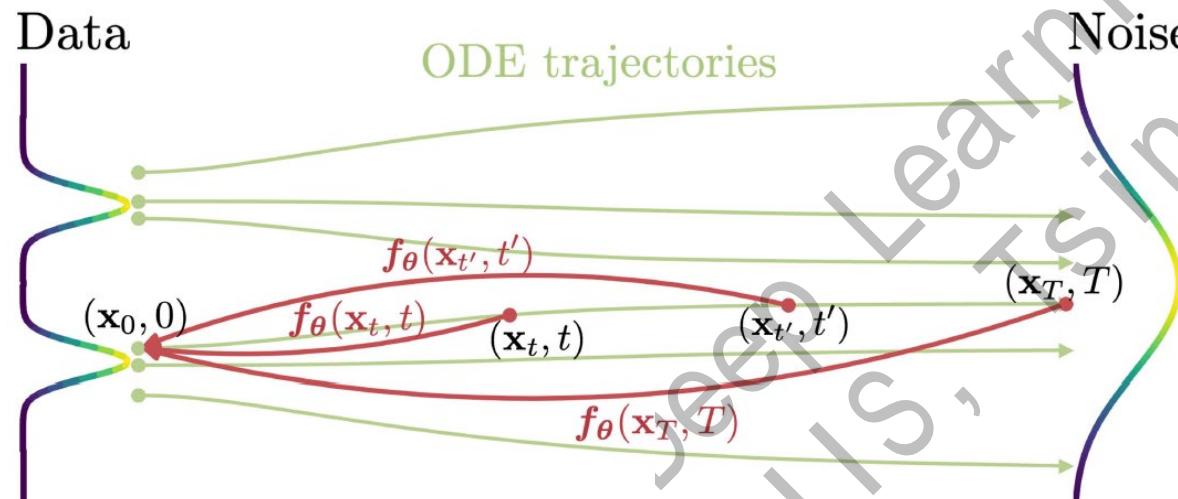
Model	FID↓	IS↑
<b>Conditional</b>		
BigGAN (Brock et al., 2018)	14.73	9.22
StyleGAN2-ADA (Karras et al., 2020a)	<b>2.42</b>	<b>10.14</b>
<b>Unconditional</b>		
StyleGAN2-ADA (Karras et al., 2020a)	2.92	9.83
NCSN (Song & Ermon, 2019)	25.32	8.87 $\pm$ .12
NCSNv2 (Song & Ermon, 2020)	10.87	8.40 $\pm$ .07
DDPM (Ho et al., 2020)	3.17	9.46 $\pm$ .11
DDPM++		
DDPM++ cont. (VP)	2.78	9.64
DDPM++ cont. (sub-VP)	2.55	9.58
DDPM++ cont. (deep, VP)	2.61	9.56
DDPM++ cont. (deep, sub-VP)	2.41	9.68
NCSN++	2.41	9.57
NCSN++ cont. (VE)	2.45	9.73
NCSN++ cont. (deep, VE)	<b>2.38</b>	<b>9.83</b>
NCSN++ cont. (deep, VE)	<b>2.20</b>	<b>9.89</b>

# High-Fidelity Generation for 1024x1024 Images



# Can we further accelerate inference?

- ODE solver still requires iterative computation
  - Can we make it faster?
- Consistency Models: use a neural network for ODE prediction!
  - Given a smooth ODE, learn  $f_\theta(x_t, t) \rightarrow x_0$  to map to trajectory origin.
  - The mapping over the same trajectory should be consistent
  - Training by distillation (more details can be found in the CM paper)



# Can we further accelerate inference?

- **Consistency Models:** use a neural network for ODE prediction!
  - Given a smooth ODE, learn  $f_\theta(x_t, t) \rightarrow x_0$  to map to trajectory origin.

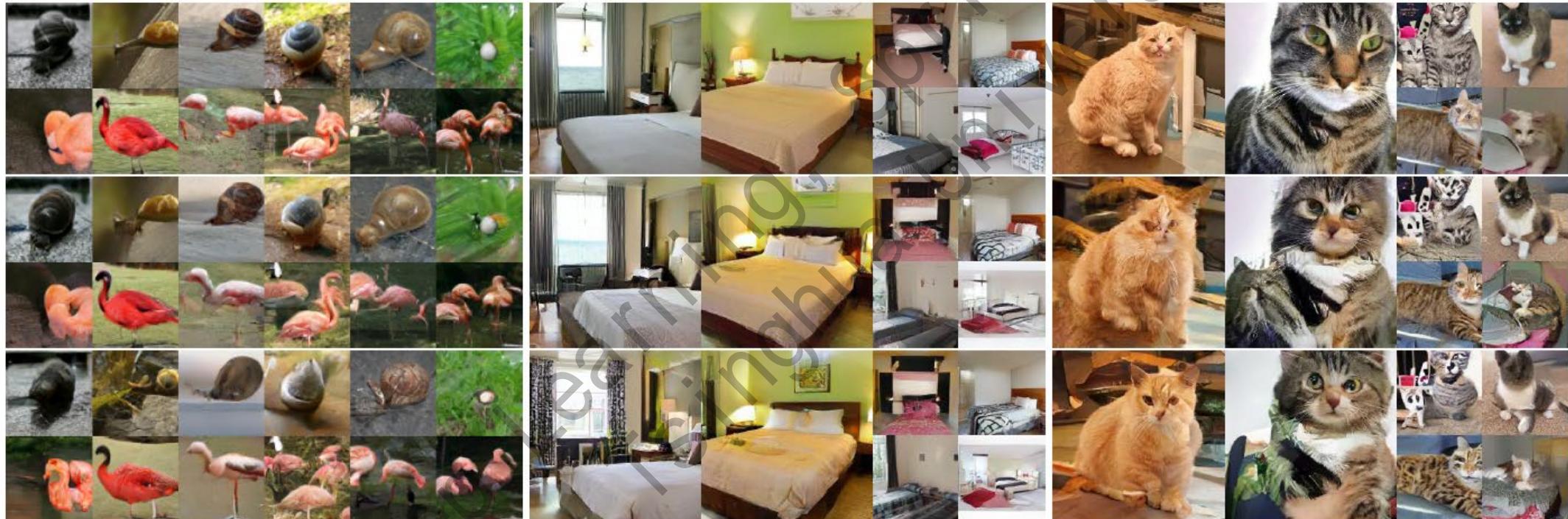
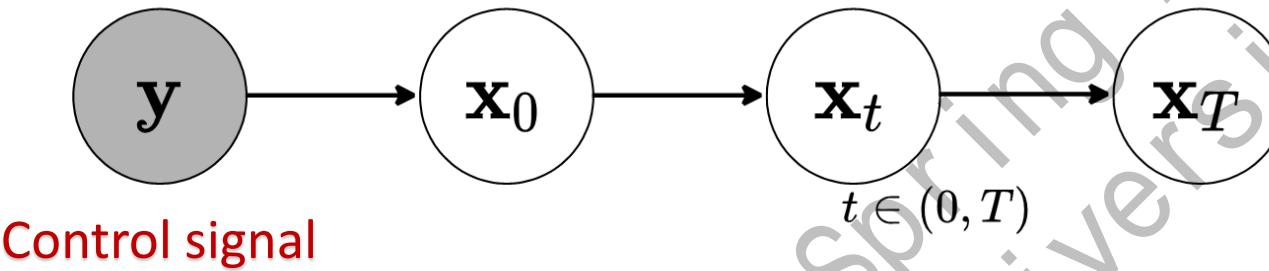


Figure 5: Samples generated by EDM (*top*), CT + single-step generation (*middle*), and CT + 2-step generation (*Bottom*). All corresponding images are generated from the same initial noise.

# Controllable Generation



- Conditional reverse-time SDE via **unconditional scores**

$$d\mathbf{x} = -\sigma^2(t) \boxed{\nabla_{\mathbf{x}} \log p_t(\mathbf{x} \mid \mathbf{y})} dt + \sigma(t) d\bar{\mathbf{w}}$$

$$d\mathbf{x} = -\sigma^2(t) [\boxed{\nabla_{\mathbf{x}} \log p_t(\mathbf{x})} + \boxed{\nabla_{\mathbf{x}} \log p_t(\mathbf{y} \mid \mathbf{x})}] dt + \sigma(t) d\bar{\mathbf{w}}$$

unconditional score,  
Trained w/o y

Trained separately or  
specified with domain knowledge

# Controllable Generation: class-conditional generation

- $y$  is the **class label**
- $p_t(y | x)$  is a time-dependent classifier



# Controllable Generation: inpainting

- **y** is the **masked image**
- $p_t(\mathbf{y} \mid \mathbf{x})$  can be approximated without training

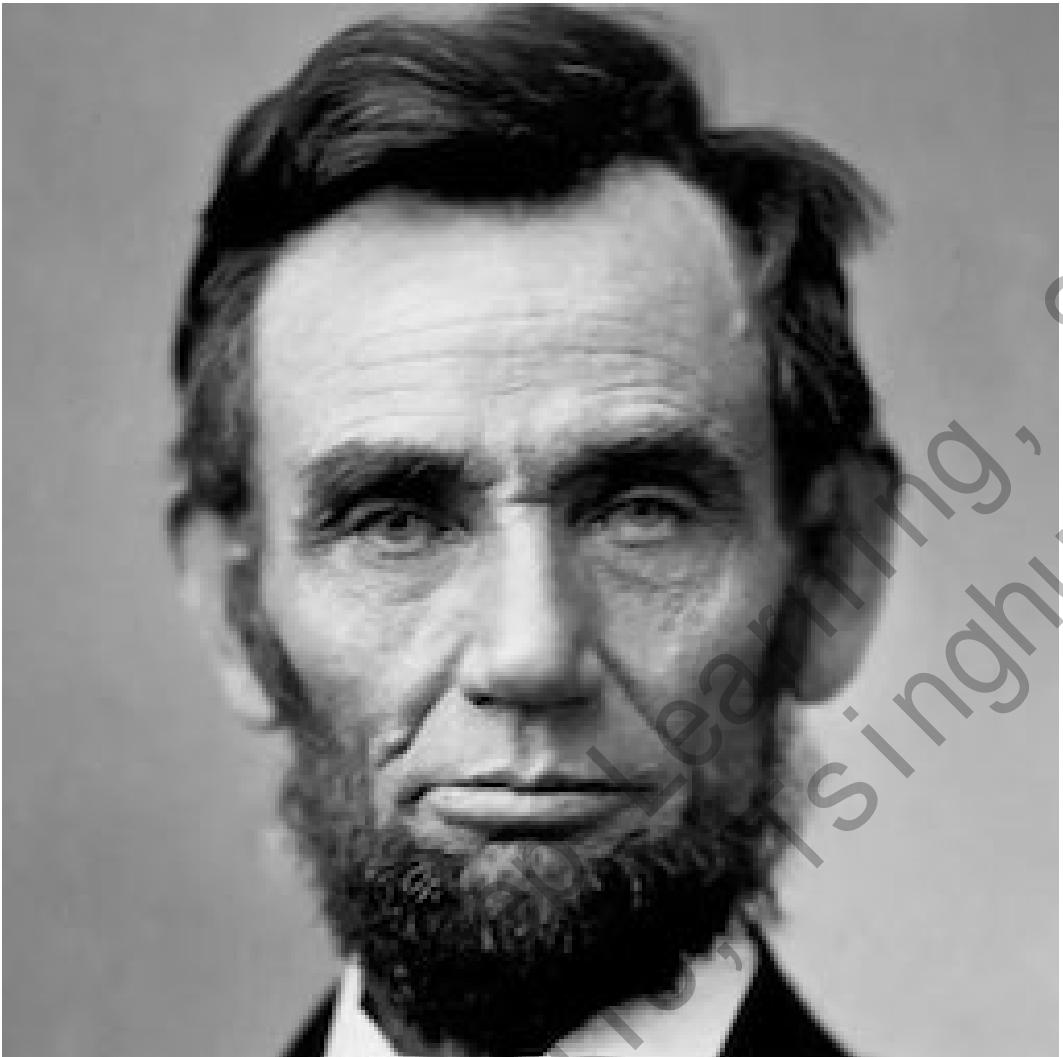


# Controllable Generation: colorization

- **y** is the gray-scale image
- $p_t(\mathbf{y} \mid \mathbf{x})$  can be approximated without training



# Controllable Generation: colorization



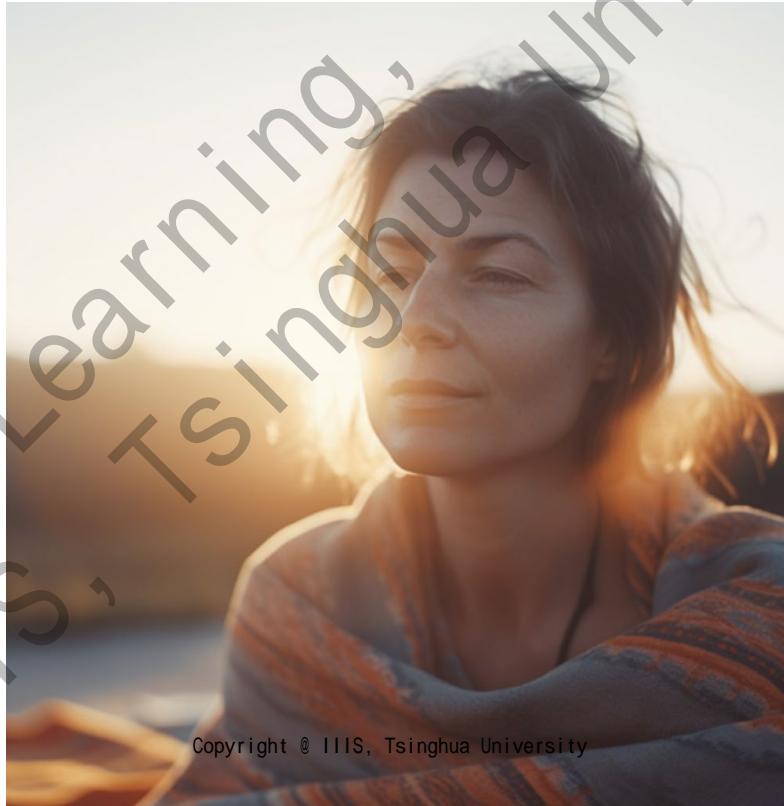
Resolution: 1024x1024

# Controllable generation: Text-guided generation

- An astronaut riding a horse in photorealistic style (Dall-E 2)



- A very attractive and natural woman, sitting on a yoga mat, breathing, eye closed, no make up, intense satisfaction, she looks like she is intensely relaxed, yoga class, sunrise, 35mm, F1: 4 (Midjourney v5)



- Cozy Scandinavian living room, there is a cat sleeping on the couch, depth of field (Midjourney v5)



# Build Your Diffusion Model

- Hugging face will be your best friend!
  - Hugging face is a platform for sharing 😊 models
  - Example: Stable Diffusion Model <https://huggingface.co/runwayml/stable-diffusion-v1-5>
- How to develop your own model?
  - For example, I want to build a model to generating 二次元 images
  - Re-training can be expensive!

We wanted to know how much time (and money) it would cost to train a Stable Diffusion model from scratch using our Streaming datasets, Composer, and MosaicML platform. Our results: it would take us 79,000 A100-hours in 13 days, for a total training cost of less than \$160,000. 2023年1月24日

Number of A100s	Throughput (images / second)	Days to Train on MosaicML Cloud	A100-hours	Approx. Cost on MosaicML Cloud
8	128.2	216.83	49,696	\$99,000
16	254.0	108.63	50,968	\$100,000
32	485.7	68.33	52,470	\$105,000
64	932.2	36.38	55,875	\$110,000
128	1864.4	18.5	62,987	\$126,000
256*	2569.4	12.83	76,735	\$160,000



mosaicml.com

<https://www.mosaicml.com/blog/training-stable-diffusion-from-scratch-costs-less-than-160k/>

# Build Your Diffusion Model

- LORA: Low-rank adaptation of large language models (MSR, 2021)
  - Low-rank hypothesis
    - Neural network models with over-parameterization reside in a low intrinsic dimension
    - Fine-tuning can be also performed with a “low-rank” fashion
  - Low-rank decomposition of additive weights
    - Model weights are frozen
    - $A$  is initialized to small Gaussian noise
    - $B$  is initialized to zero
    - <https://github.com/microsoft/LoRA>
  - Lora becomes extremely popular these days ...
    - Some examples on next slide ...

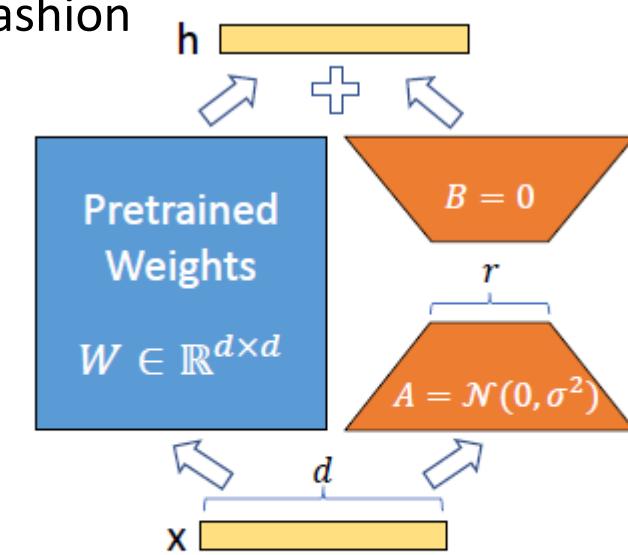
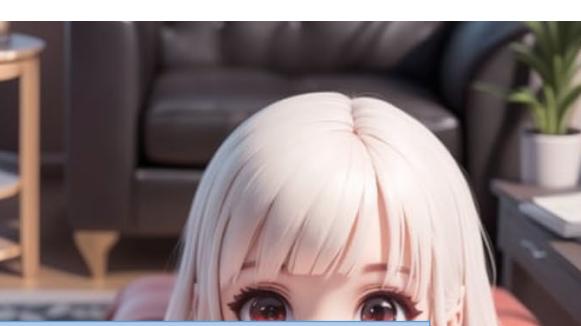
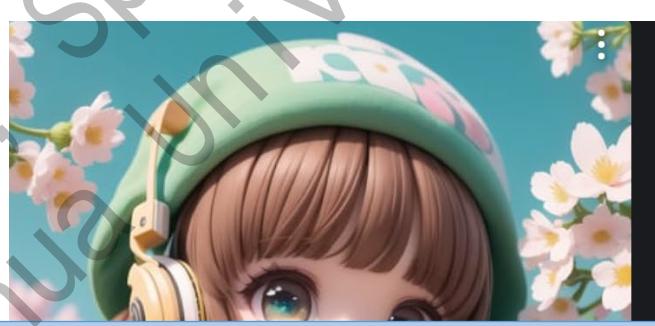
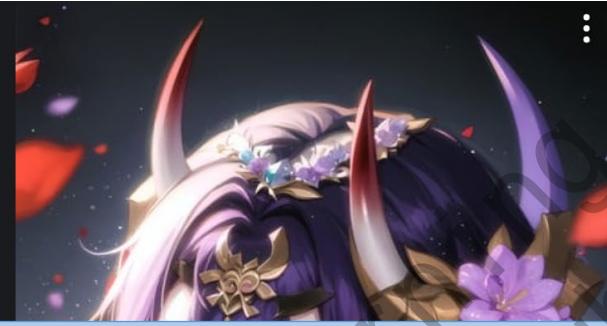
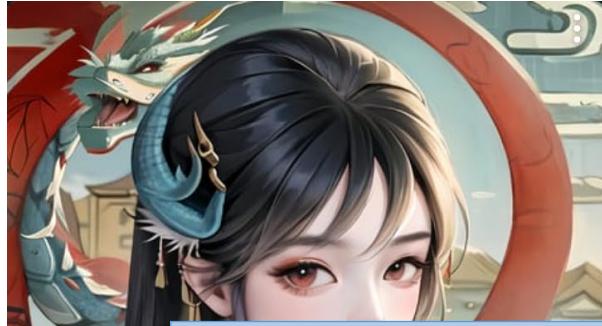


Figure 1: Our reparametrization. We only train  $A$  and  $B$ .

# Build Your Diffusion Model: Examples

- Fashion Girl



**But what if I want a fine-grained control beyond text?**



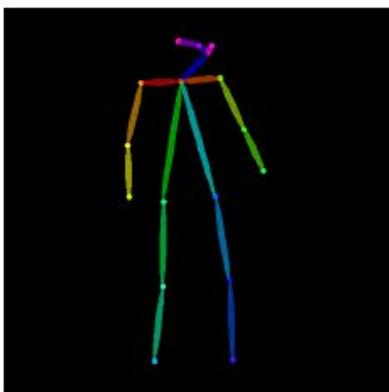
# Control Your Model

- ControlNet: Adding Conditional Control to Text-to-Image Diffusion Models (Stanford, 2023)
  - Let the model also condition on additional signal
    - Example: Canny edge detection
    - <https://github.com/Illyasviel/ControlNet>



# Control Your Model

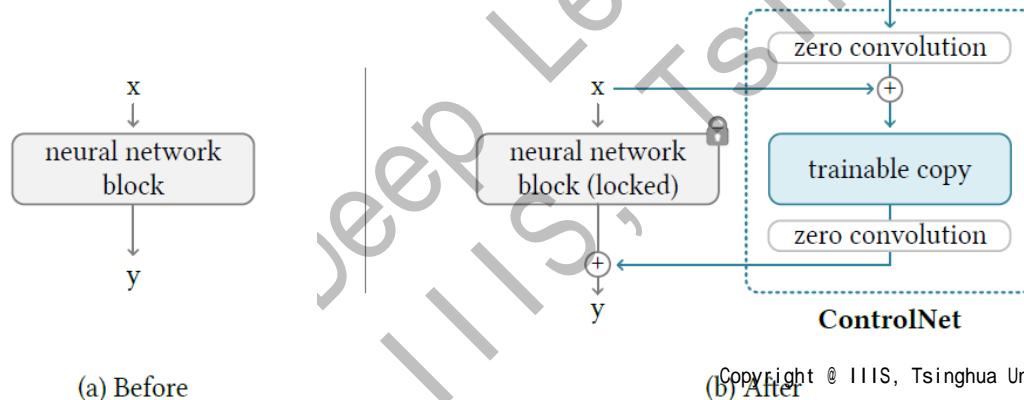
- ControlNet: Adding Conditional Control to Text-to-Image Diffusion Models (Stanford, 2023)
  - Let the model also condition on additional signal
    - Example: Human Motion
    - <https://github.com/llyasviel/ControlNet>



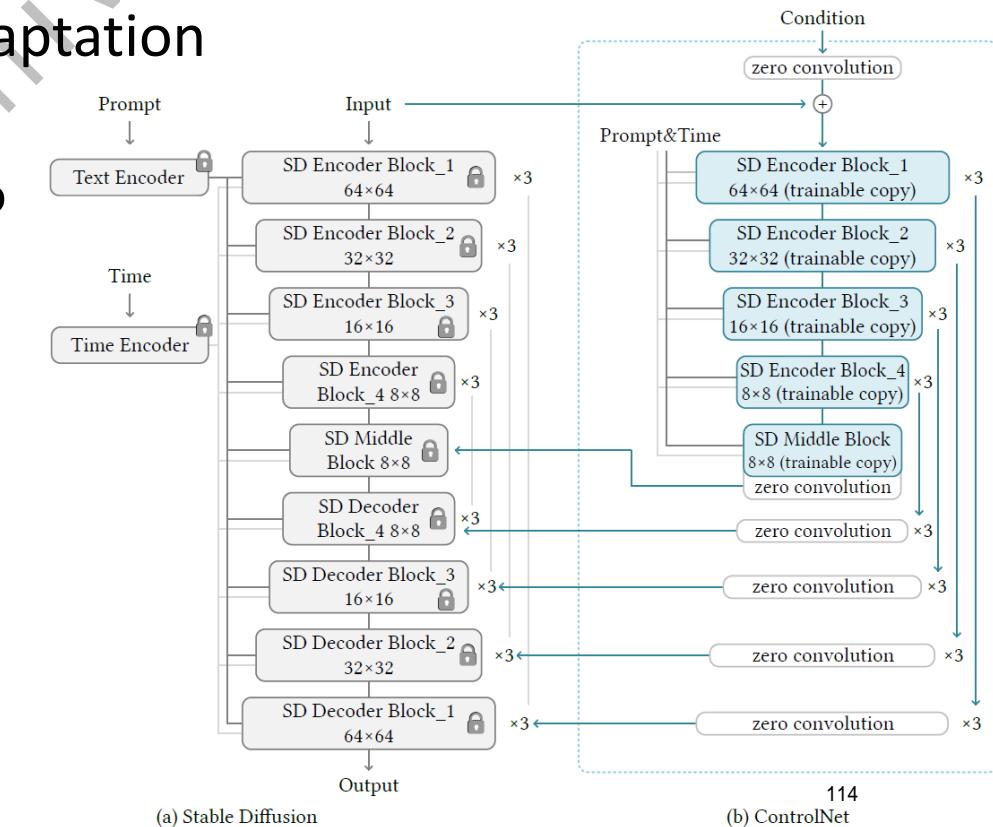
“astronaut”

# Control Your Model

- ControlNet: Adding Conditional Control to Text-to-Image Diffusion Models (Stanford, 2023)
  - Similar idea to Lora: frozen weight + small adaptation
  - Key techniques:
    - Zero-convolution: 1x1 conv-layer initializes to zero
    - Trainable copy for adaptation initialization
    - Repeated additive for conditioning



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- More ControlNet examples
  - Home designer



# Control Your Model

- ControlNet: Adding Conditional Control to Diffusion Models (Stanford, 2023)
- More ControlNet examples
  - Home designer
  - Even video!
    - more in future lectures...



# Summary: Diffusion Model

- Advanced Generative Model
  - Diffusion Model and Score-Based Model
    - Score matching for gradients of log probability
  - Training & Inference (with condition)
    - Noise conditioned network
    - Langevin dynamics and SDE for fast sampling
    - Conditioned generation without the need of retraining
- Frontier AIGC: LORA and ControlNet
  - Hugging Face is your friend

# Generative Model (Summary)

- Goal of generative model
  - Learn a distribution  $p(x)$  to generate samples and unsupervised learning
- Models so far
  - Energy-based model
    - $p(x) = \frac{1}{Z} \exp(-E(x))$ , powerful representation but hard to sample
  - Variational auto-encoder
    - $p(x, z) = p(z)p(x|z)$ , variational inference as an approximate method
  - Generative adversarial net
    - $G(z)$  and  $D(x)$ , an implicit model with high generation quality and unstable training
  - Normalizing flow
    - $x = f(z)$ , best mathematical properties but the most restricted representation
  - Score-based models
    - $s(x) = \nabla_x p(x; \theta)$ , highest generation quality + stable training, but generation is slow

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**Coming next: generating data samples beyond images!**

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# Thanks!

Deep Learning, Spring 2025  
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