

Trustworthy AI Systems

-- Generative Modeling (Part II)

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Last Lecture

- Generative Adversarial Network
 - DCGAN
 - Conditional GAN
 - CycleGAN

- Neural Style Transfer

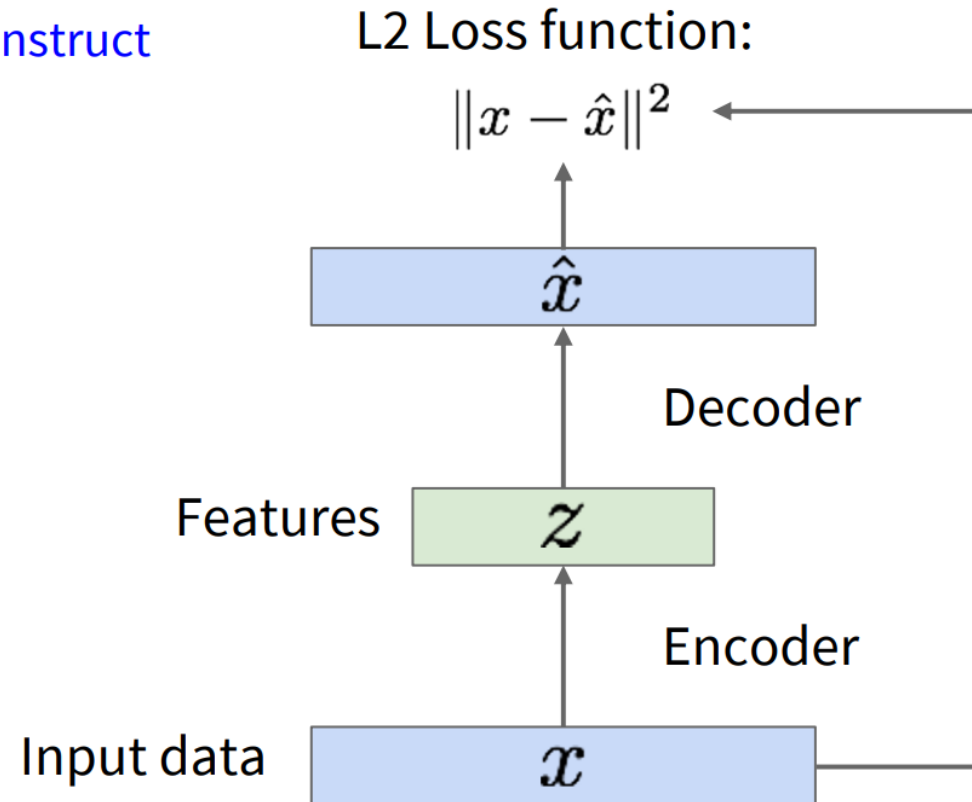
This Lecture

- Variational Autoencoders
- Diffusion Models

Autoencoder

Train such that features can be used to reconstruct original data

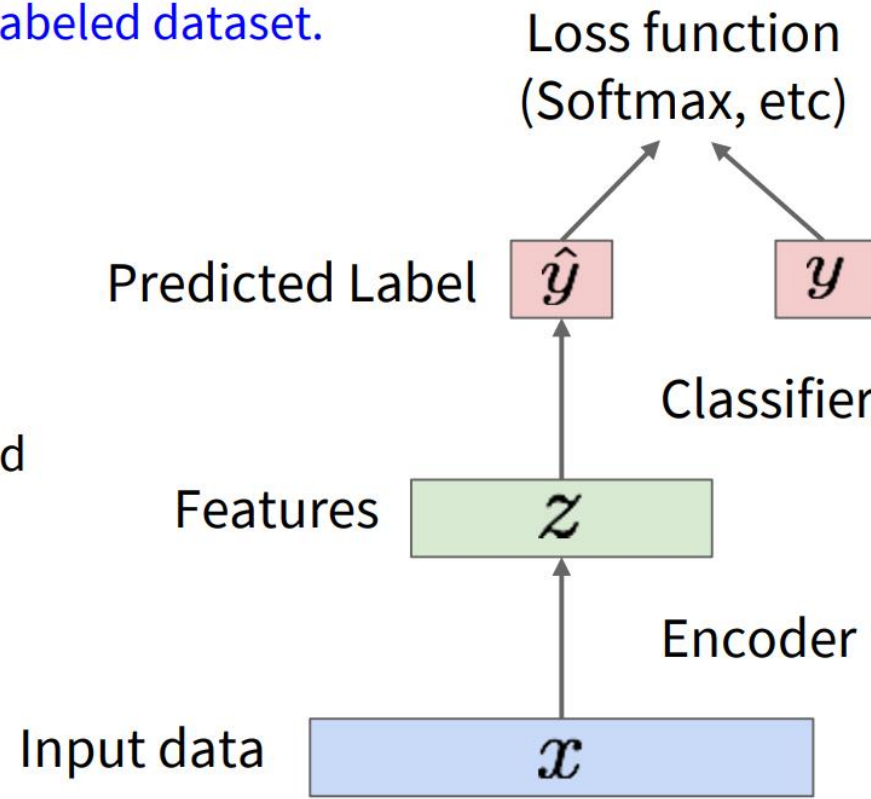
Doesn't use labels!



Application of Autoencoder

Transfer from large, unlabeled dataset to small, labeled dataset.

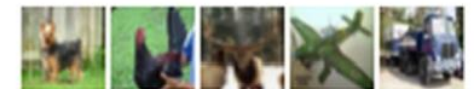
Encoder can be used to initialize a supervised model



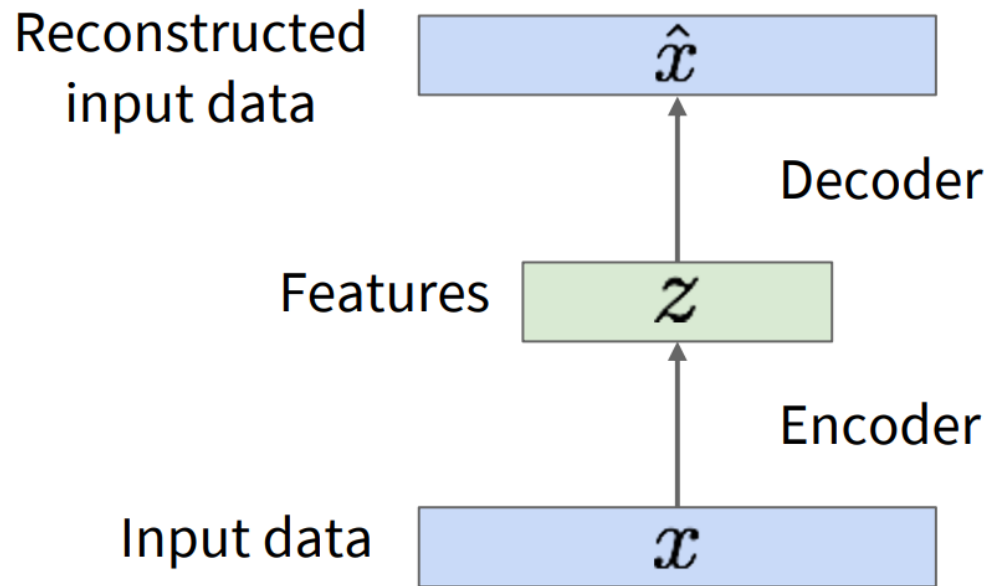
Fine-tune encoder jointly with classifier

bird plane
dog deer truck

Train for final task (sometimes with small data)



The Limitation of Autoencoder



Autoencoders can reconstruct data, and can learn features to initialize a supervised model

Features capture factors of variation in training data.

But we can't generate new images from an autoencoder because we don't know the space of z .

How do we make autoencoder a generative model?

Variational Autoencoders

Probabilistic spin on autoencoders - will let us sample from the model to generate data!

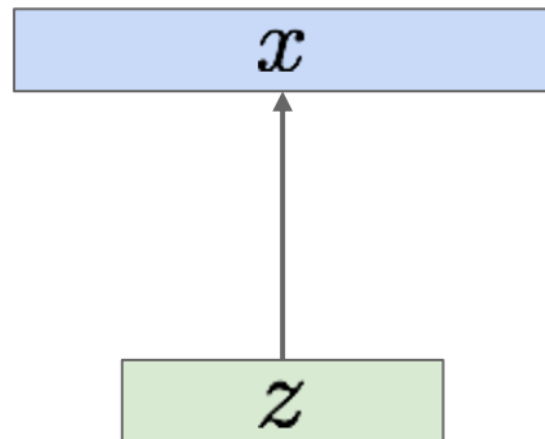
Assume training data $\{x^{(i)}\}_{i=1}^N$ is generated from the distribution of unobserved (latent) representation z

Sample from true conditional

$$p_{\theta^*}(x | z^{(i)})$$

Sample from true prior

$$z^{(i)} \sim p_{\theta^*}(z)$$



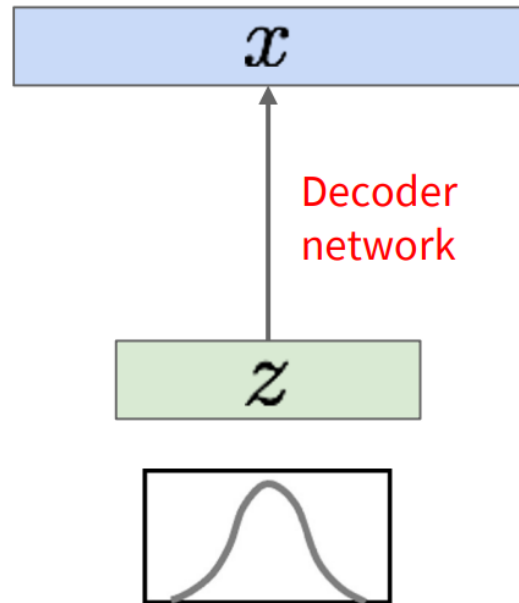
Intuition (remember from autoencoders!): x is an image, z is latent factors used to generate x : attributes, orientation, etc.

Variational Autoencoders



Sample from
true conditional
 $p_{\theta^*}(x | z^{(i)})$

Sample from
true prior
 $z^{(i)} \sim p_{\theta^*}(z)$



We want to estimate the true parameters θ^* of this generative model given training data x .

How should we represent this model?

Choose prior $p(z)$ to be simple, e.g. Gaussian.
Reasonable for latent attributes, e.g. pose, how much smile.

Conditional $p(x|z)$ is complex (generates image)
 \Rightarrow represent with neural network

Variational Autoencoders: Intractability

Data likelihood: $p_{\theta}(x) = \int p_{\theta}(z)p_{\theta}(x|z)dz$

 Intractable to compute $p(x|z)$ for every z !

$$\log p(x) \approx \log \frac{1}{k} \sum_{i=1}^k p(x|z^{(i)}), \text{ where } z^{(i)} \sim p(z)$$

Monte Carlo estimation is too high variance

KL measures how one probability distribution P is different from a second, Q .

$$D_{\text{KL}}(P \parallel Q) = \sum_{x \in \mathcal{X}} P(x) \log \left(\frac{P(x)}{Q(x)} \right).$$

Variational Autoencoders

Math trick: Taking expectation wrt. z (using encoder network)

$$\log p_{\theta}(x^{(i)}) = \mathbf{E}_{z \sim q_{\phi}(z|x^{(i)})} \left[\log p_{\theta}(x^{(i)}) \right] \quad (p_{\theta}(x^{(i)})) \text{ Does not depend on } z$$

We want to maximize the data likelihood

$$= \mathbf{E}_z \left[\log \frac{p_{\theta}(x^{(i)} | z) p_{\theta}(z)}{p_{\theta}(z | x^{(i)})} \right] \quad (\text{Bayes' Rule})$$

$$= \mathbf{E}_z \left[\log \frac{p_{\theta}(x^{(i)} | z) p_{\theta}(z) q_{\phi}(z | x^{(i)})}{p_{\theta}(z | x^{(i)}) q_{\phi}(z | x^{(i)})} \right] \quad (\text{Multiply by constant})$$

$$= \mathbf{E}_z \left[\log p_{\theta}(x^{(i)} | z) \right] - \mathbf{E}_z \left[\log \frac{q_{\phi}(z | x^{(i)})}{p_{\theta}(z)} \right] + \mathbf{E}_z \left[\log \frac{q_{\phi}(z | x^{(i)})}{p_{\theta}(z | x^{(i)})} \right] \quad (\text{Logarithms})$$

$$= \mathbf{E}_z \left[\log p_{\theta}(x^{(i)} | z) \right] - D_{\text{KL}}(q_{\phi}(z | x^{(i)}) \parallel p_{\theta}(z)) + D_{\text{KL}}(q_{\phi}(z | x^{(i)}) \parallel p_{\theta}(z | x^{(i)}))$$

↑
Decoder network gives $p_{\theta}(x|z)$, can compute estimate of this term through sampling (need some trick to differentiate through sampling).

↑
This KL term (between Gaussians for encoder and z prior) has nice closed-form solution!

↑
 $p_{\theta}(z|x)$ intractable (saw earlier), can't compute this KL term :(But we know KL divergence always ≥ 0 .

Variational Autoencoders

$$\log p_{\theta}(x^{(i)}) = \mathbf{E}_{z \sim q_{\phi}(z|x^{(i)})} \left[\log p_{\theta}(x^{(i)}) \right] \quad (p_{\theta}(x^{(i)} \text{ Does not depend on } z)$$

$$= \mathbf{E}_z \left[\log \frac{p_{\theta}(x^{(i)} | z) p_{\theta}(z)}{p_{\theta}(z | x^{(i)})} \right] \quad (\text{Bayes' Rule})$$

Decoder:
reconstruct
the input data

$$= \mathbf{E}_z \left[\log \frac{p_{\theta}(x^{(i)} | z) p_{\theta}(z) q_{\phi}(z | x^{(i)})}{p_{\theta}(z | x^{(i)}) q_{\phi}(z | x^{(i)})} \right] \quad (\text{Multiply by constant})$$

Encoder:
make approximate
posterior distribution
close to prior

$$= \mathbf{E}_z \left[\log p_{\theta}(x^{(i)} | z) \right] - \mathbf{E}_z \left[\log \frac{q_{\phi}(z | x^{(i)})}{p_{\theta}(z)} \right] + \mathbf{E}_z \left[\log \frac{q_{\phi}(z | x^{(i)})}{p_{\theta}(z | x^{(i)})} \right] \quad (\text{Logarithms})$$

$$= \underbrace{\mathbf{E}_z \left[\log p_{\theta}(x^{(i)} | z) \right] - D_{KL}(q_{\phi}(z | x^{(i)}) || p_{\theta}(z))}_{\mathcal{L}(x^{(i)}, \theta, \phi)} + \underbrace{D_{KL}(q_{\phi}(z | x^{(i)}) || p_{\theta}(z | x^{(i)}))}_{\geq 0}$$

Tractable lower bound which we can take
gradient of and optimize! ($p_{\theta}(x|z)$ differentiable,
KL term differentiable)

Reparameterization in VAE

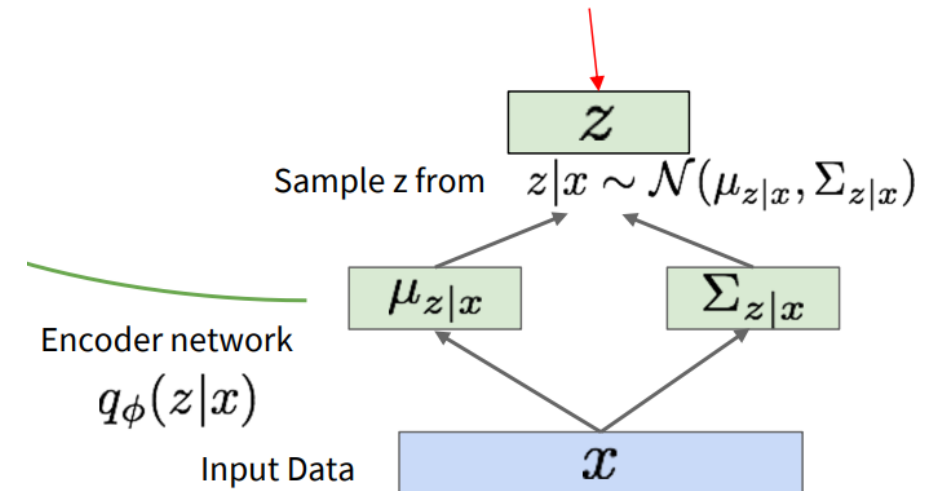
- Generate **NEW**: Sampling is required to model the probabilistic nature of latent space.
- This sampling operation introduce stochasticity and therefore cannot be differentiated.
- Backpropagation relies on computing gradients of deterministic (i.e., non-random) operations.

Reparameterization trick to make sampling differentiable:

$$\text{Sample } \epsilon \sim \mathcal{N}(0, I)$$

$$z = \mu_{z|x} + \epsilon \sigma_{z|x}$$

Not part of the computation graph!



Variational Autoencoders

$$\underbrace{\mathbf{E}_z \left[\log p_\theta(x^{(i)} | z) \right]}_{\mathcal{L}(x^{(i)}, \theta, \phi)} - \underbrace{D_{KL}(q_\phi(z | x^{(i)}) || p_\theta(z))}_{\text{blue box}}$$

$$D_{KL}(\mathcal{N}(\mu_{z|x}, \Sigma_{z|x}) || \mathcal{N}(0, I))$$

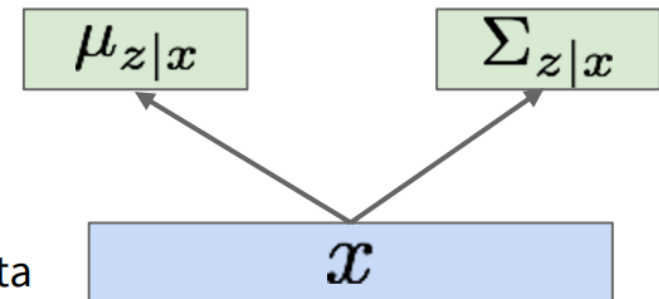
Have analytical solution

Make approximate posterior distribution close to prior

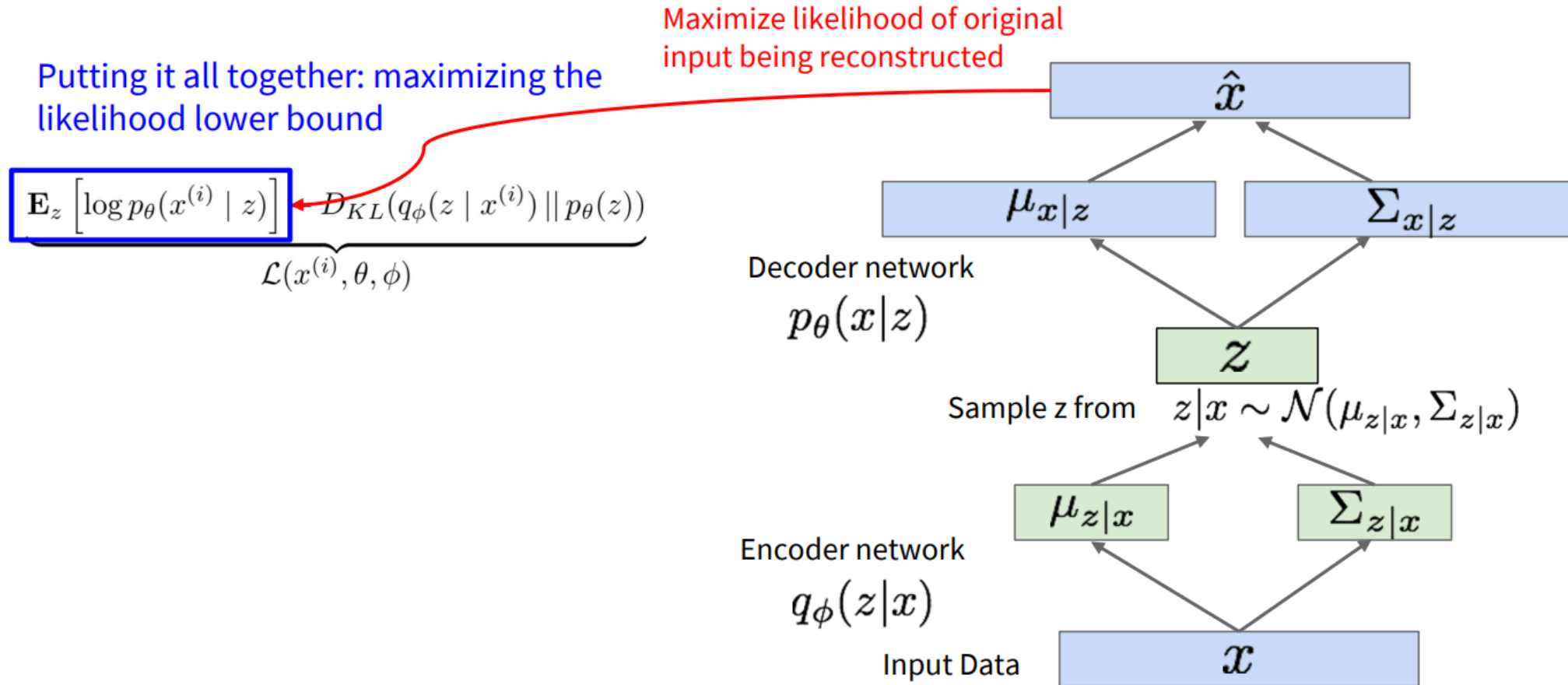
Encoder network

$$q_\phi(z|x)$$

Input Data



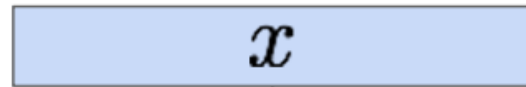
Variational Autoencoders



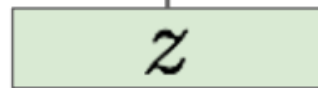
Variational Autoencoders: Generating Data

Our assumption about data generation process

Sample from true conditional
 $p_{\theta^*}(x | z^{(i)})$

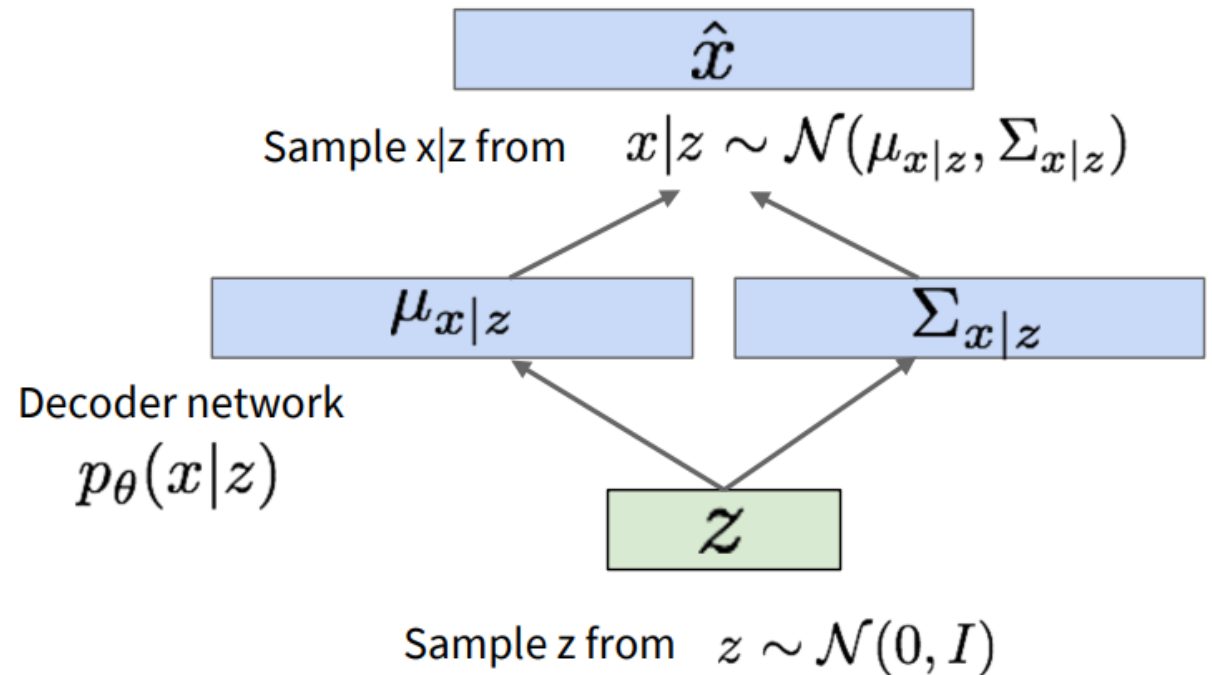


Sample from true prior
 $z^{(i)} \sim p_{\theta^*}(z)$



Decoder network

Now given a trained VAE:
use decoder network & sample z from prior!



Variational Autoencoders: Generating Data!



32x32 CIFAR-10



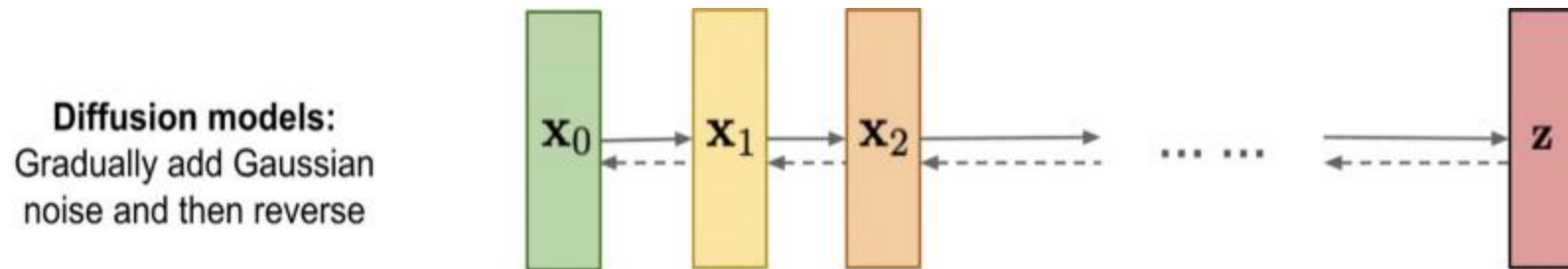
Labeled Faces in the Wild

Take a Break

Diffusion

Idea: Estimating and analyzing small step sizes is more tractable/easier than a single step from random noise to the learned distribution

Convert a well-known and simple base distribution (like a Gaussian) to the target (data) distribution iteratively, with small step sizes, via a Markov chain



Forward Diffusion Process

- Noise added can be parameterized by:

$$q(\mathbf{x}_t|\mathbf{x}_{t-1}) = \mathcal{N}(\mathbf{x}_t; \sqrt{1 - \beta_t}\mathbf{x}_{t-1}, \beta_t\mathbf{I}) \quad q(\mathbf{x}_{1:T}|\mathbf{x}_0) = \prod_{t=1}^T q(\mathbf{x}_t|\mathbf{x}_{t-1}) \quad \{\beta_t \in (0, 1)\}_{t=1}^T$$

Vary the parameters of the Gaussian according to a *noise schedule*


- You can prove with some math that as T approaches infinity, you eventually end up with an Isotropic Gaussian (i.e. pure random noise)
- Note: forward process is fixed

Forward Diffusion Process

Reparameterization trick:

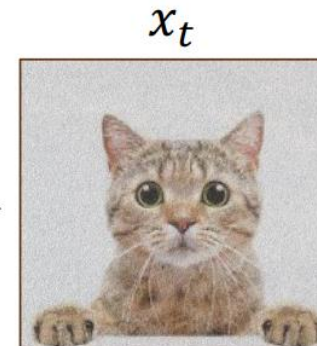
$$q(\mathbf{x}_t | \mathbf{x}_0) = \mathcal{N}\left(\sqrt{\bar{\alpha}_t} \mathbf{x}_0, (1 - \bar{\alpha}_t) \mathbf{I}\right)$$

$$\begin{aligned} \alpha_t &= 1 - \beta_t \\ \bar{\alpha}_t &= \prod_{i=1}^t \alpha_i \end{aligned}$$

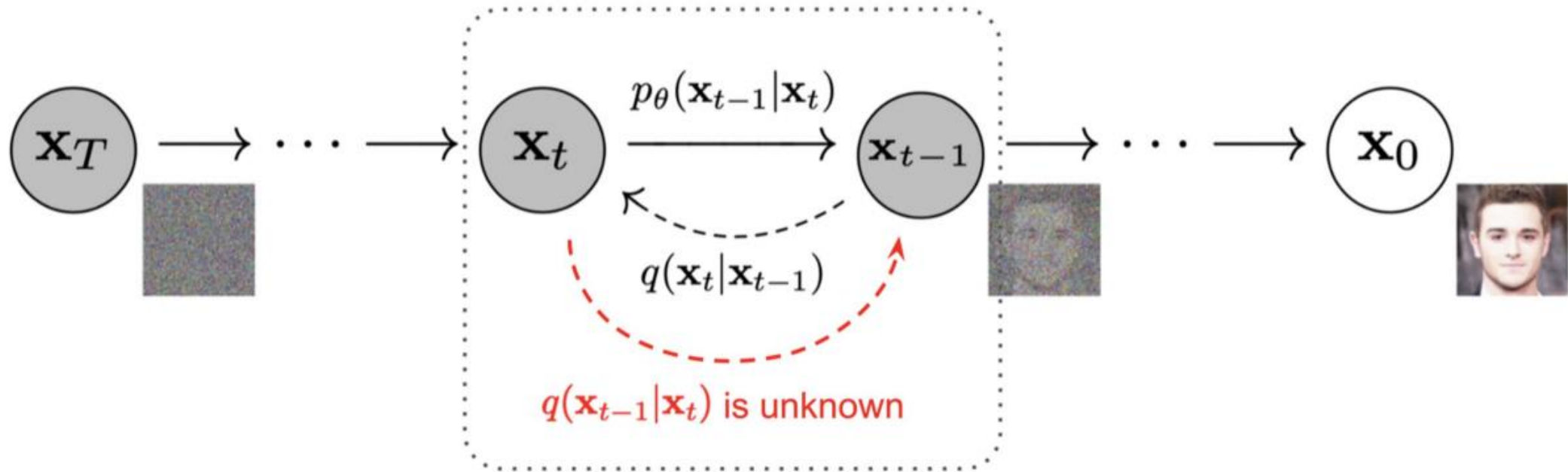
1. Sample an image from the dataset: 
2. Sample noise $\epsilon \sim N(0, \mathbf{I})$ (from a **standard** normal distribution)
3. Scale the image by $\sqrt{\alpha_t}$: $\sqrt{\alpha_t} x_0$

where $\alpha_t = 1 - \beta_t$
 $\bar{\alpha}_t = \prod_{i=1}^t \alpha_i$

4. Add $\sqrt{1 - \bar{\alpha}_t} \epsilon$: $\sqrt{\bar{\alpha}_t} x_0 + \sqrt{1 - \bar{\alpha}_t} \epsilon$ \longrightarrow



Reverse Diffusion Process



The goal of a diffusion model is to learn the reverse denoising process to iteratively undo the forward process

Distribution in Reverse Process

Turns out that for small enough forward steps, i.e. $\{\beta_t \in (0, 1)\}_{t=1}^T$

the reverse process step $q(\mathbf{x}_{t-1} | \mathbf{x}_t)$ can be estimated as a Gaussian distribution too

Therefore, we can parametrize the *learned* reverse process as

$$p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t) = \mathcal{N}(\mathbf{x}_{t-1}; \boldsymbol{\mu}_{\theta}(\mathbf{x}_t, t), \boldsymbol{\Sigma}_{\theta}(\mathbf{x}_t, t))$$

In practice, Σ is just the identity matrix, so we only need to learn the mean of the distribution

Loss Function

$$\begin{aligned}
 \mathbf{x}_t &= \sqrt{\alpha_t} \mathbf{x}_{t-1} + \sqrt{1 - \alpha_t} \boldsymbol{\epsilon}_{t-1} \\
 &= \sqrt{\alpha_t \alpha_{t-1}} \mathbf{x}_{t-2} + \sqrt{1 - \alpha_t \alpha_{t-1}} \bar{\boldsymbol{\epsilon}}_{t-2} \\
 &= \dots \\
 &= \sqrt{\bar{\alpha}_t} \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \boldsymbol{\epsilon}
 \end{aligned}$$

$$\begin{aligned}
 L_t &= \mathbb{E}_{\mathbf{x}_0, \boldsymbol{\epsilon}} \left[\frac{1}{2 \|\boldsymbol{\Sigma}_\theta(\mathbf{x}_t, t)\|_2^2} \|\tilde{\boldsymbol{\mu}}_t(\mathbf{x}_t, \mathbf{x}_0) - \boldsymbol{\mu}_\theta(\mathbf{x}_t, t)\|_2^2 \right] \\
 &= \mathbb{E}_{\mathbf{x}_0, \boldsymbol{\epsilon}} \left[\frac{1}{2 \|\boldsymbol{\Sigma}_\theta\|_2^2} \left\| \frac{1}{\sqrt{\alpha_t}} \left(\mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1 - \bar{\alpha}_t}} \boldsymbol{\epsilon}_t \right) - \frac{1}{\sqrt{\alpha_t}} \left(\mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1 - \bar{\alpha}_t}} \boldsymbol{\epsilon}_\theta(\mathbf{x}_t, t) \right) \right\|_2^2 \right] \\
 &= \mathbb{E}_{\mathbf{x}_0, \boldsymbol{\epsilon}} \left[\frac{(1 - \alpha_t)^2}{2 \alpha_t (1 - \bar{\alpha}_t) \|\boldsymbol{\Sigma}_\theta\|_2^2} \|\boldsymbol{\epsilon}_t - \boldsymbol{\epsilon}_\theta(\mathbf{x}_t, t)\|_2^2 \right] \\
 &= \mathbb{E}_{\mathbf{x}_0, \boldsymbol{\epsilon}} \left[\frac{(1 - \alpha_t)^2}{2 \alpha_t (1 - \bar{\alpha}_t) \|\boldsymbol{\Sigma}_\theta\|_2^2} \|\boldsymbol{\epsilon}_t - \boldsymbol{\epsilon}_\theta(\sqrt{\bar{\alpha}_t} \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \boldsymbol{\epsilon}_t, t)\|_2^2 \right]
 \end{aligned}$$

Training and Sampling

Algorithm 1 Training

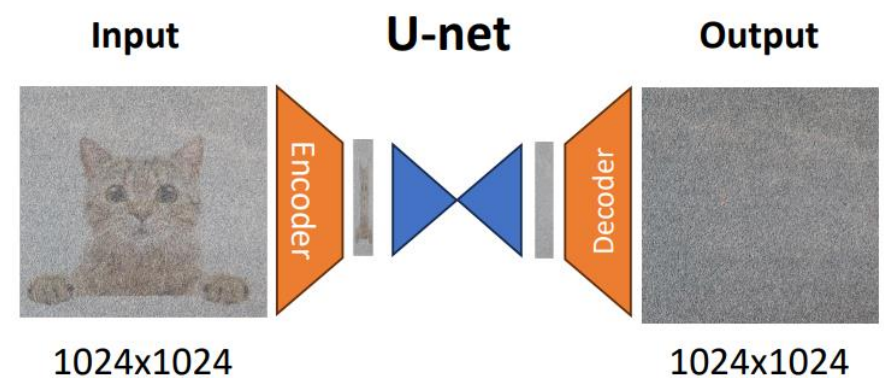
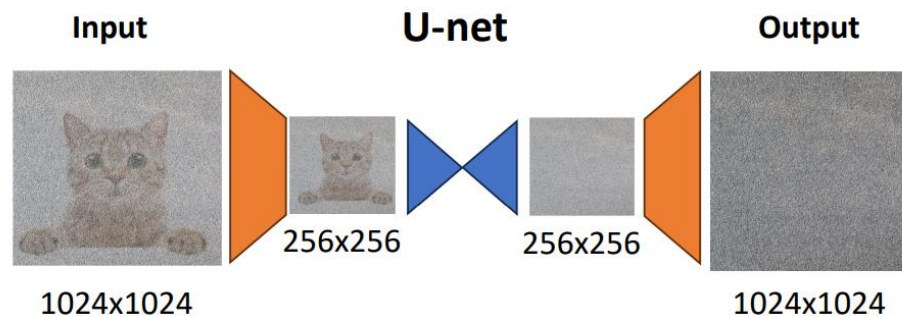
- 1: **repeat**
 - 2: $\mathbf{x}_0 \sim q(\mathbf{x}_0)$
 - 3: $t \sim \text{Uniform}(\{1, \dots, T\})$
 - 4: $\boldsymbol{\epsilon} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$
 - 5: Take gradient descent step on
$$\nabla_{\theta} \left\| \boldsymbol{\epsilon} - \boldsymbol{\epsilon}_{\theta}(\sqrt{\bar{\alpha}_t}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

Algorithm 2 Sampling

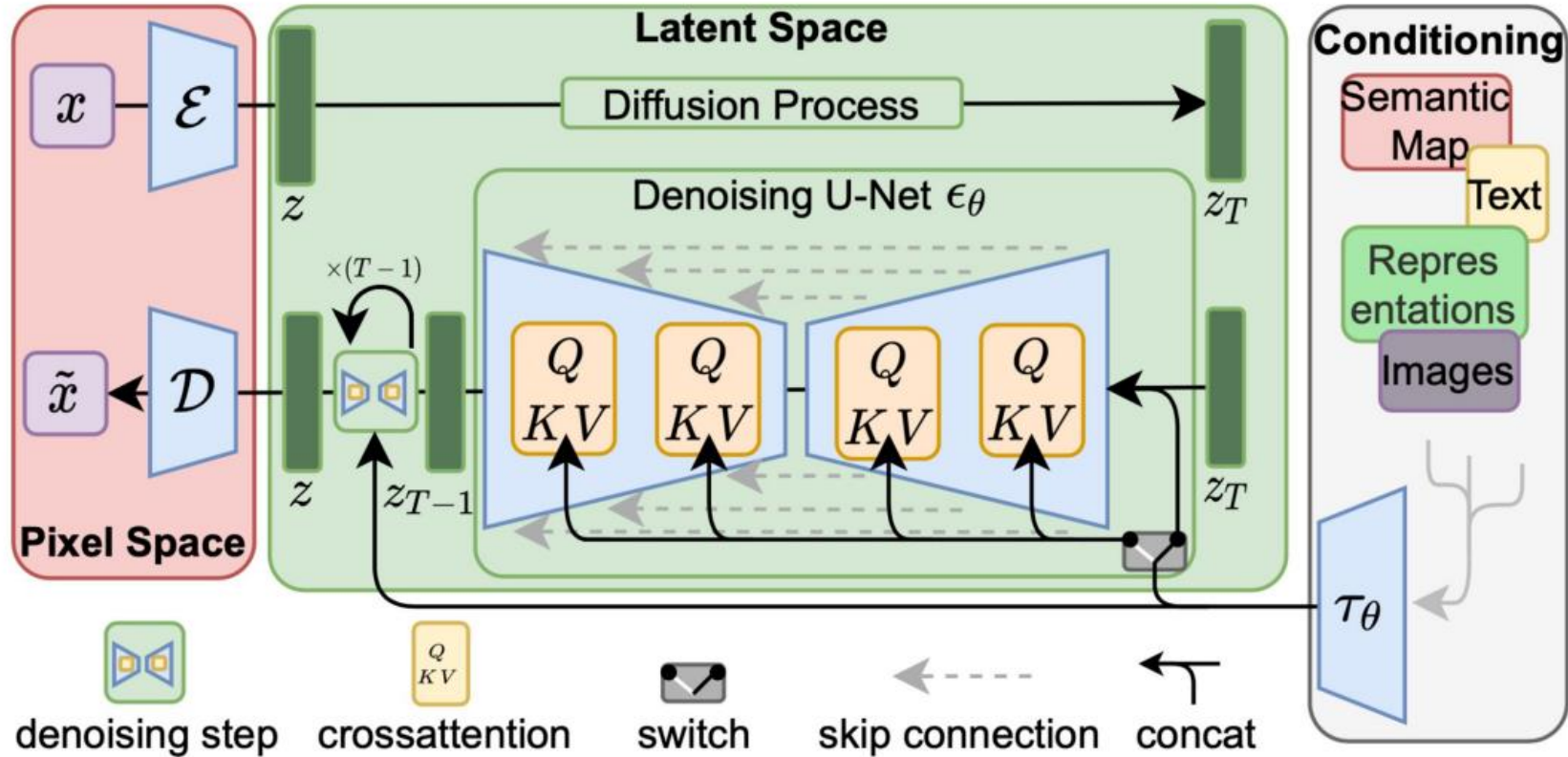
- 1: $\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$
 - 2: **for** $t = T, \dots, 1$ **do**
 - 3: $\mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ if $t > 1$, else $\mathbf{z} = \mathbf{0}$
 - 4: $\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left(\mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1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

U-Net Problem

- Operating in the input space is very computationally expensive
 - Generate Low-Resolution + Upsample
 - Generate in Latent Space



Stable Diffusion



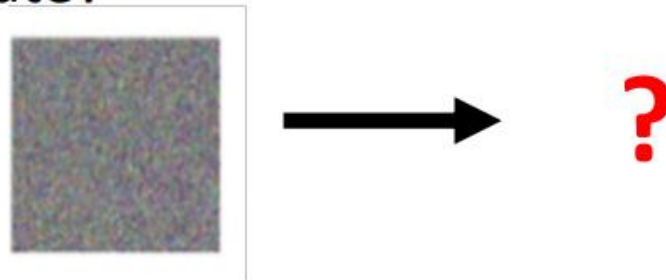
High-Resolution Image Synthesis with Latent Diffusion Models (CVPR 2022)
<https://ommer-lab.com/research/latent-diffusion-models/>

Guided/Conditioned Diffusion

Lets say we train a diffusion model on images of cats and dogs:



If we start from random noise, and generate a new image, what will the model generate?

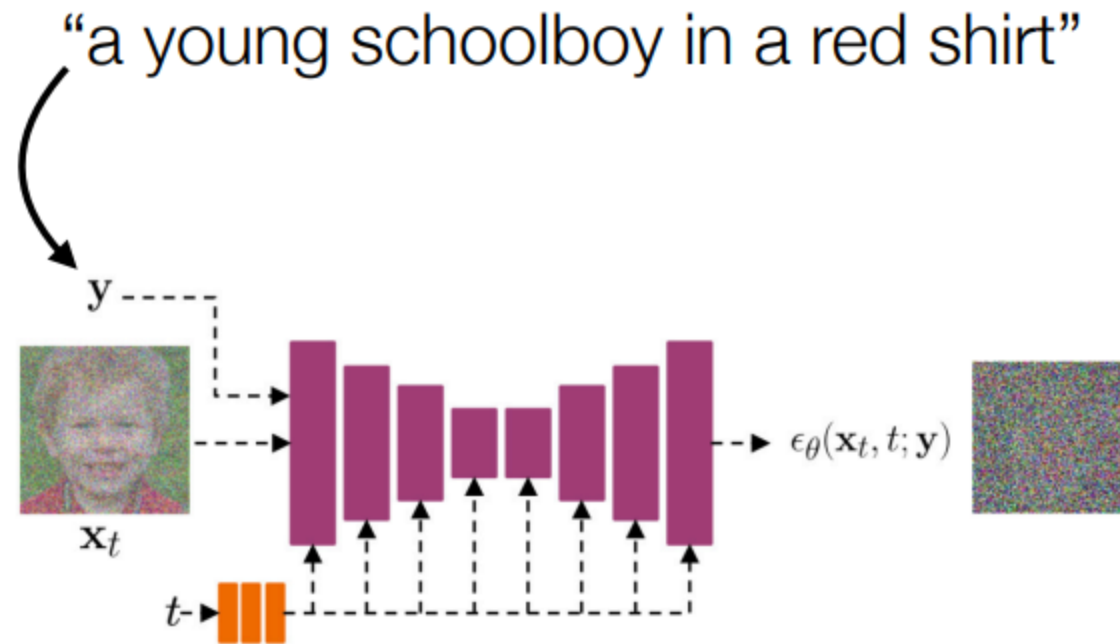


How to Control Diffusion Models?

- Explicit conditioning
- Classifier guidance
- Classifier-free guidance

Explicit conditioning

Use an Image-Text dataset



Classifier Guidance

Given a Gaussian distribution of \mathbf{x}

$$\nabla_{\mathbf{x}} \log p(\mathbf{x}) = \nabla_{\mathbf{x}} \left(-\frac{1}{2\sigma^2} (\mathbf{x} - \boldsymbol{\mu})^2 \right) = -\frac{\mathbf{x} - \boldsymbol{\mu}}{\sigma^2} = -\frac{\boldsymbol{\epsilon}}{\sigma} \text{ where } \boldsymbol{\epsilon} \sim \mathcal{N}(\mathbf{0}, \mathbf{I}). \text{ Recall that}$$

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

$$\mathbf{s}_{\theta}(\mathbf{x}_t, t) \approx \nabla_{\mathbf{x}_t} \log q(\mathbf{x}_t) = \mathbb{E}_{q(\mathbf{x}_0)} [\nabla_{\mathbf{x}_t} q(\mathbf{x}_t | \mathbf{x}_0)] = \mathbb{E}_{q(\mathbf{x}_0)} \left[-\frac{\boldsymbol{\epsilon}_{\theta}(\mathbf{x}_t, t)}{\sqrt{1 - \bar{\alpha}_t}} \right] = -\frac{\boldsymbol{\epsilon}_{\theta}(\mathbf{x}_t, t)}{\sqrt{1 - \bar{\alpha}_t}}$$

Score function for the joint distribution:

$$\begin{aligned} \nabla_{\mathbf{x}_t} \log q(\mathbf{x}_t, y) &= \nabla_{\mathbf{x}_t} \log q(\mathbf{x}_t) + \nabla_{\mathbf{x}_t} \log q(y | \mathbf{x}_t) \\ &\approx -\frac{1}{\sqrt{1 - \bar{\alpha}_t}} \boldsymbol{\epsilon}_{\theta}(\mathbf{x}_t, t) + \nabla_{\mathbf{x}_t} \log f_{\phi}(y | \mathbf{x}_t) \\ &= -\frac{1}{\sqrt{1 - \bar{\alpha}_t}} (\boldsymbol{\epsilon}_{\theta}(\mathbf{x}_t, t) - \sqrt{1 - \bar{\alpha}_t} \nabla_{\mathbf{x}_t} \log f_{\phi}(y | \mathbf{x}_t)) \end{aligned}$$

Classifier Guidance

Thus, a new classifier-guided predictor $\bar{\epsilon}_\theta$ would take the form as following,

$$\bar{\epsilon}_\theta(\mathbf{x}_t, t) = \epsilon_\theta(x_t, t) - \sqrt{1 - \bar{\alpha}_t} \nabla_{\mathbf{x}_t} \log f_\phi(y|\mathbf{x}_t)$$

To control the strength of the classifier guidance, we can add a weight w to the delta part,

$$\bar{\epsilon}_\theta(\mathbf{x}_t, t) = \epsilon_\theta(x_t, t) - \sqrt{1 - \bar{\alpha}_t} w \nabla_{\mathbf{x}_t} \log f_\phi(y|\mathbf{x}_t)$$

Algorithm 1 Classifier guided diffusion sampling, given a diffusion model $(\mu_\theta(x_t), \Sigma_\theta(x_t))$, classifier $f_\phi(y|x_t)$, and gradient scale s .

Input: class label y , gradient scale s
 $x_T \leftarrow$ sample from $\mathcal{N}(0, \mathbf{I})$
for all t from T to 1 **do**
 $\mu, \Sigma \leftarrow \mu_\theta(x_t), \Sigma_\theta(x_t)$
 $x_{t-1} \leftarrow$ sample from $\mathcal{N}(\mu + s\Sigma \nabla_{x_t} \log f_\phi(y|x_t), \Sigma)$
end for
return x_0

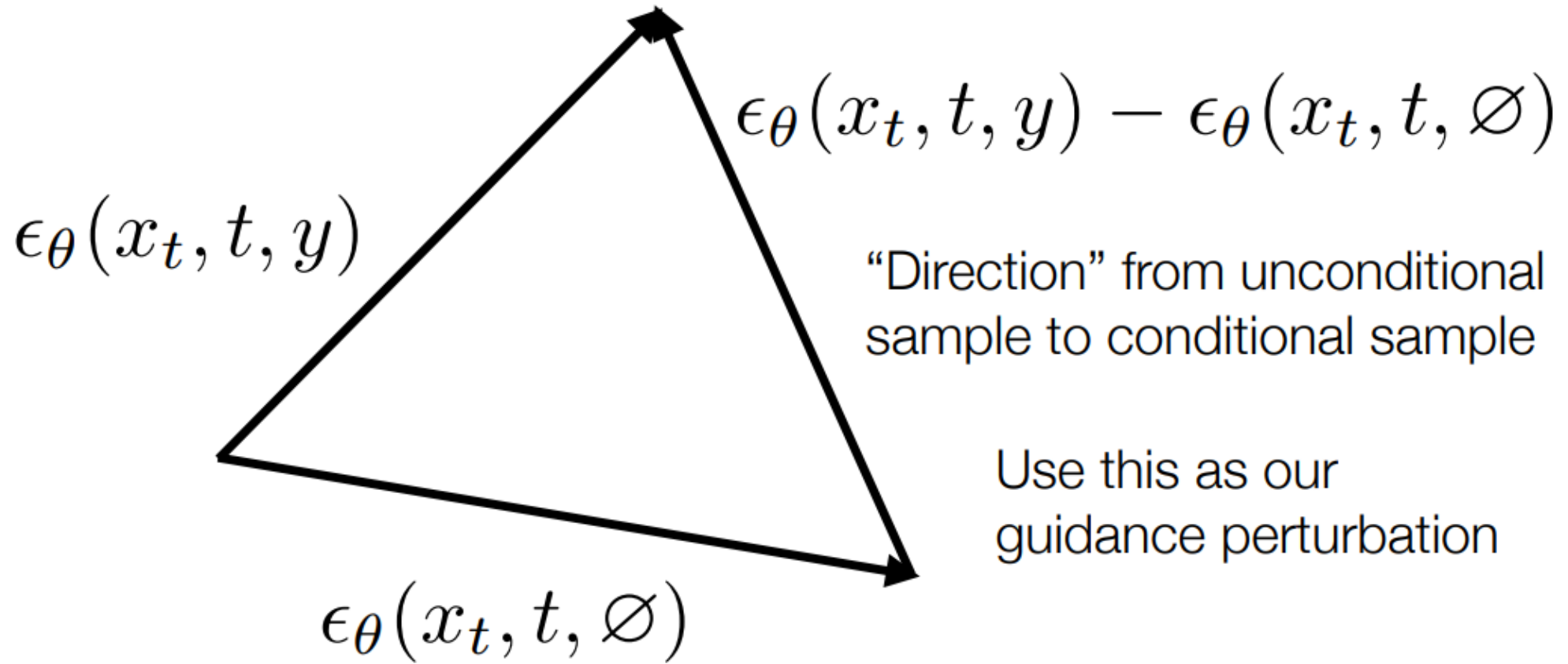
Problems with Classifier Guidance

- Need to fine-tune or re-train a classifier on noisy data
- Need a pre-trained classification model
 - What if we want to use any text prompt as input?

Classifier Free Guidance

Idea: Use the diffusion model itself to get perturbations for guidance

Classifier Free Guidance



Classifier Free Guidance

- A conditional diffusion model is trained on pair data (x, y) , where the conditioning information y get discarded at random

$$\begin{aligned}\nabla_{\mathbf{x}_t} \log p(y|\mathbf{x}_t) &= \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t|y) - \nabla_{\mathbf{x}_t} \log p(\mathbf{x}_t) \\ &= -\frac{1}{\sqrt{1 - \bar{\alpha}_t}} \left(\epsilon_{\theta}(\mathbf{x}_t, t, y) - \epsilon_{\theta}(\mathbf{x}_t, t) \right)\end{aligned}$$

Classifier-guided modified score

$$\begin{aligned}\bar{\epsilon}_{\theta}(\mathbf{x}_t, t, y) &= \epsilon_{\theta}(\mathbf{x}_t, t, y) - \sqrt{1 - \bar{\alpha}_t} w \nabla_{\mathbf{x}_t} \log p(y|\mathbf{x}_t) \\ &= \epsilon_{\theta}(\mathbf{x}_t, t, y) + w \left(\epsilon_{\theta}(\mathbf{x}_t, t, y) - \epsilon_{\theta}(\mathbf{x}_t, t) \right) \\ &= (w + 1) \epsilon_{\theta}(\mathbf{x}_t, t, y) - w \epsilon_{\theta}(\mathbf{x}_t, t)\end{aligned}$$

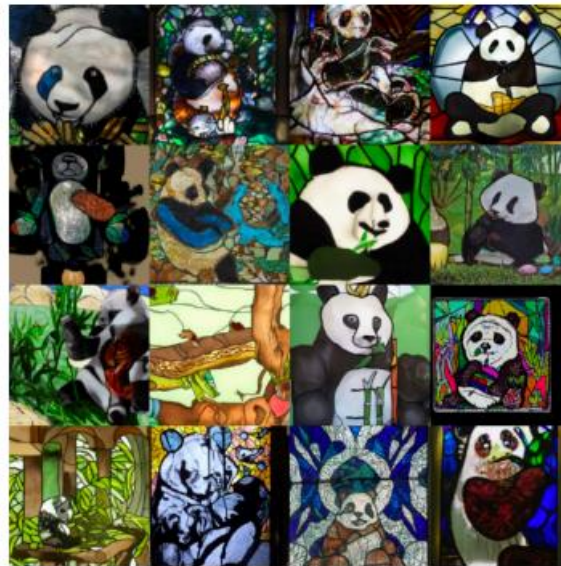
Classifier Free Guidance: Text2Image

Our new noise estimate will then be:

$$\tilde{\epsilon}(x_t, t, y) = \epsilon_{\theta}(x_t, t, \emptyset) + \underbrace{\gamma(\epsilon_{\theta}(x_t, t, y) - \epsilon_{\theta}(x_t, t, \emptyset))}_{\text{"Direction" from unconditional to conditional}}$$

“Direction” from unconditional to conditional

“A stained glass window of a panda eating bamboo”



$\gamma = 1$



$\gamma = 3$

References

- <https://www.eecs.umich.edu/courses/eecs442-ahowens/fa23/slides/lec11-diffusion.pdf>
- <https://lilianweng.github.io/posts/2021-07-11-diffusion-models/>
- https://cs231n.stanford.edu/slides/2024/lecture_13.pdf