## Generative Models Discriminative vs Generative

**Understanding Machine Learning** 

## Archetype: Naïve Bayes

#### probabilistic model:

$$P(Y|X_1, \cdots, X_p) = \frac{P(Y, X_1, \cdots, X_p)}{P(X_1, \cdots, X_p)} = \frac{P(Y)P(X_1, \cdots, X_p|Y)}{P(X_1, \cdots, X_p)} \propto P(Y)P(X_1, \cdots, X_p|Y)$$
Bayes' rule constant to be estimated

#### approach:

- 1. estimate  $P(Y, X) \rightarrow$  generative model (can be used to generate new samples)
- 2. calculate P(Y|X) from  $P(Y,X) \rightarrow$  used for discriminative task (classification)

## Independence Assumption

(naïve) assumption: conditional independence of features given target

$$P(X_j|Y,X_1,\cdots,X_{j-1},X_{j+1},\cdots,X_p) = P(X_j|Y)$$

$$\Rightarrow P(Y|X_1,\dots,X_p) = \frac{P(Y)\prod_{j=1}^p P(X_j|Y)}{P(X_1,\dots,X_p)}$$

- > independent feature contributions (ignoring feature correlations)
- > robust against curse of dimensionality

#### Estimation of Feature Contributions

separate estimations of  $P(X_j|Y)$  for each feature

requires assumption of distributions (e.g., Gaussian naïve Bayes) or non-parametric methods (kernel density estimation)

Gaussian feature likelihoods:

$$P(x_{ij}|y) = \frac{1}{\sqrt{2\pi\sigma_{y,j}^2}} \exp\left(-\frac{(x_{ij}-\mu_{y,j})^2}{2\sigma_{y,j}^2}\right)$$

parameter estimation (e.g., mean and variance of Gaussians) can be done with maximum likelihood method (y known in training)

→ no Bayesian methods needed

#### Maximum a Posteriori Classification

$$\hat{y}_i = \underset{y}{\operatorname{argmax}} P(y) \prod_{j=1}^p P(x_{ij}|y)$$

despite potentially inaccurate probability estimates (due to naïve independence assumption), good identification of correct class via maximum probability

→ bad for regression tasks (if independence assumption is too naïve, i.e., features are correlated)

## Generative vs Discriminative Models

generative models: predict joint probability P(Y, X) (what allows to create new data samples) or directly generate new data samples

or just  $P(X) \rightarrow$  unsupervised (or self-supervised) learning

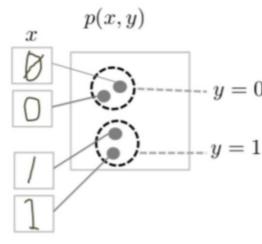
discriminative models: predict conditional probability (or probability distribution for regression) P(Y|X) or directly output (label for classification, real value for regression)

task of generative models more difficult: model full data distribution rather than merely find patterns in inputs to distinguish outputs

discriminative model



generative model



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## Naïve Bayes and Logistic Regression

generative-discriminative pair of classification algorithms

- binary case: logit of naïve Bayes' outputs,  $\log\left(\frac{P(y_i=1|x_i)}{P(y_i=0|x_i)}\right)$ , corresponds to output of logistic regression's linear predictor
- for discrete inputs or Gaussian naïve Bayes: naïve Bayes can be reparametrized as linear classifier

for discriminative task: identical in asymptotic limit (infinite training samples) if independence assumption holds (otherwise naïve Bayes less accurate)

naïve Bayes has greater bias but lower variance than logistic regression  $\rightarrow$  to be preferred for scarce training data (if bias, i.e., independence assumption, correct)

#### Data Generation

generative models can be used for discriminative tasks (although potentially inferior to direct discriminative methods)

but generative methods do more than discriminative ones: model full data distribution

→ allows generation of new data samples (can be images, text, video, audio, code like SQL or Python, proteins, materials, time series, structured data, ...)

large (auto-regressive) language models examples of generative models

#### Generative Al

Depending on the application, there are currently two dominant approaches for generative AI:

• text generation: LLMs

• image synthesis: diffusion models

## Image Synthesis

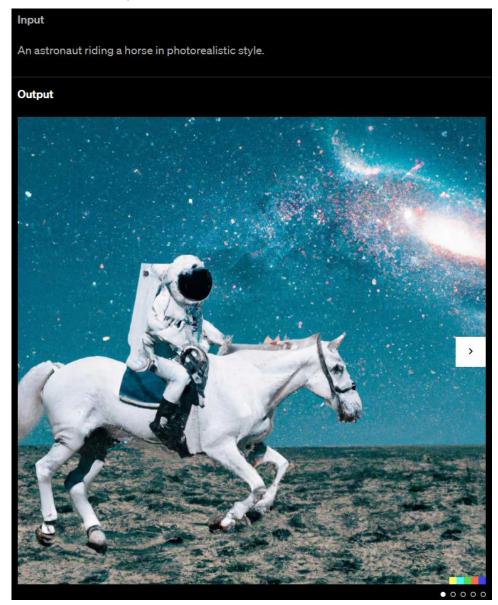
idea: generate new images as variations of training data

condition generation on text prompts: text-to-image

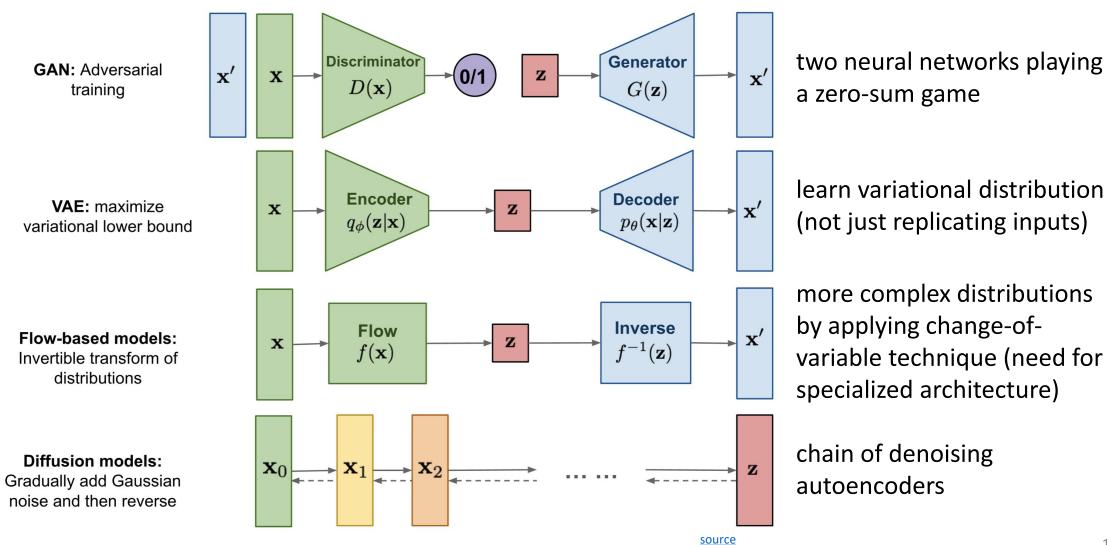
trade-off between diversity and fidelity

SOTA: (guided) diffusion models

example: DALL-E 2



## Different Types of Generative Models



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# Generative Adversarial Networks (GAN)

## Indirect Training via Discriminator

two neural networks playing a zero-sum game:

- the generator network G generating new (fake) samples
- the discriminator network D trying to distinguish between real and fake samples

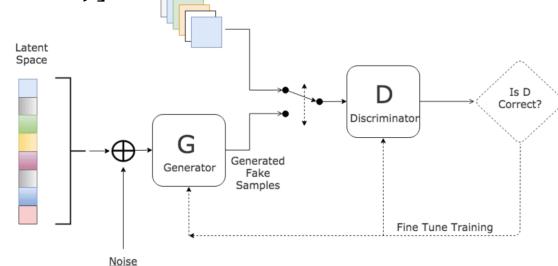
idea: G not trained directly to minimize reconstruction error of real samples, but to fool D  $\rightarrow$  self-supervised approach

#### Formulation

common loss for generator and discriminator:

$$L(\boldsymbol{x}_i) = E_{\boldsymbol{x} \sim p_r(\boldsymbol{x})}[\ln D(\boldsymbol{x}_i)] + E_{\boldsymbol{x} \sim p_g(\boldsymbol{x})}[\ln (1 - D(\boldsymbol{x}_i))]$$

- G trying to minimize
- D trying to maximize



generator: decomposition into latent space (parameters of generator network) and noise (sampled from, e.g., Gaussian distribution)

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## Properties

implicit generative model: do not estimate likelihood function

for optimal D, GAN loss quantifies similarity between generative data distribution  $p_{\it g}$  and real data distribution  $p_{\it r}$  by Jensen-Shannon divergence

$$D_{JS}(p||q) = \frac{1}{2} D_{KL}\left(p||\frac{p+q}{2}\right) + \frac{1}{2} D_{KL}\left(q||\frac{p+q}{2}\right)$$

for optimal values of both G and D:  $p_g = p_r$  and D = 0.5

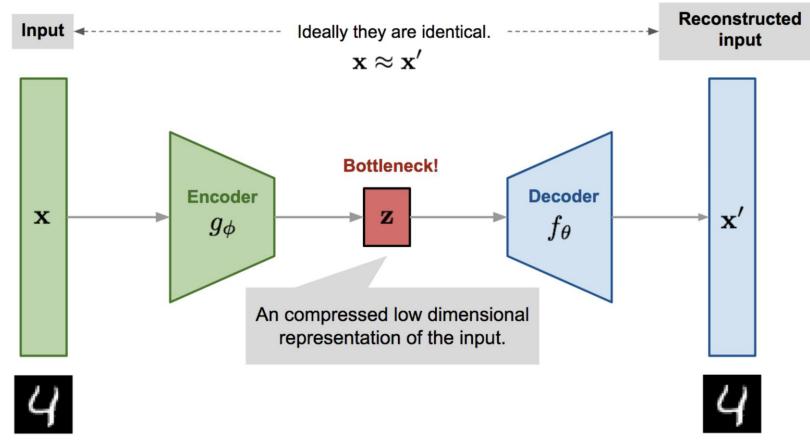
issue: potentially unstable training

## Variational Autoencoders (VAE)

## Recap: Autoencoder

(deep) encoder network
(deep) decoder network
learned together by
minimizing differences
between original input and
reconstructed input
(expressed as losses)

compressed intermediate representation: dimensionality reduction



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## Autoencoder Architecture for Generative Tasks

goal: generation of variations of input data rather than compressed representation

→ learn variational distribution instead of identity function

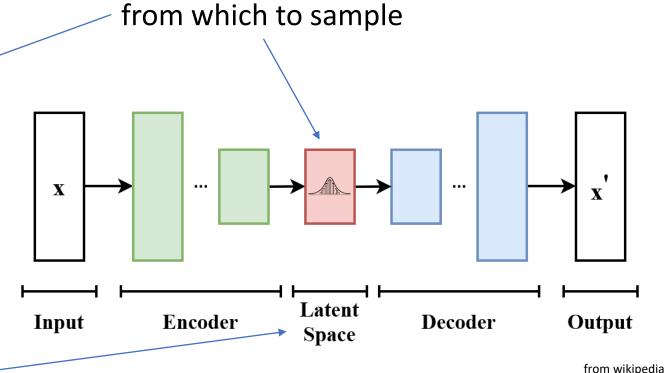
to be precise: parametrized variational distribution of latent encoding variables z

prior (simple distribution, in usual VAE:

Gaussian):  $p_{\theta}(\mathbf{z})$ 

posterior: 
$$p_{\theta}(\mathbf{z}|\mathbf{x}) = \frac{p_{\theta}(\mathbf{x}|\mathbf{z})p_{\theta}(\mathbf{z})}{\int p_{\theta}(\mathbf{x}|\mathbf{z})p_{\theta}(\mathbf{z})d\mathbf{z}}$$

 $p_{\theta}(x)$ : mixture of Gaussians



Variational Bayesian Method

## Encoder and Decoder Networks

encoder: find posterior  $p_{\theta}(\mathbf{z}|\mathbf{x})$  unfortunately, generally intractable

 $\rightarrow$  approximate by  $q_{\phi}(\mathbf{z}|\mathbf{x})$ 

VAE:  $q_{\phi}(\mathbf{z}|\mathbf{x})$  expressed by neural network with weights  $\phi$ 

 $\Rightarrow$  amortized inference:  $q_{\phi}(z|x)$  learned in training, z inferred from x in prediction (sharing variational parameters across all data points)

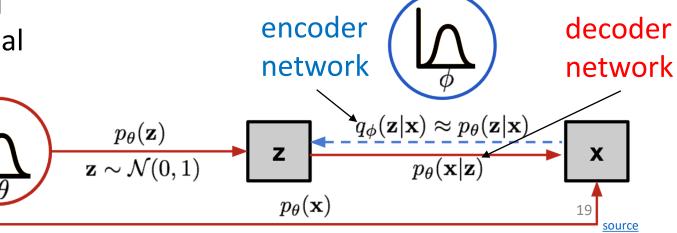
in VAE: network weights  $oldsymbol{ heta}$ 

decoder: generate new sample  $x_i$ 

- 1. sample  $z_i$  (from Gaussian)
- 2. generate  $x_i$  (similar to real data)

 $\rightarrow$  maximize:  $p_{\theta}(x_i) = \int p_{\theta}(x_i|\mathbf{z}) p_{\theta}(\mathbf{z}) d\mathbf{z}$ 

(integral over z expensive  $\rightarrow$  use only likely codes z given input x: need for encoder)



#### VAE Loss: ELBO

VAE loss function to be minimized according to network weights:

$$L(\mathbf{x}_i; \boldsymbol{\theta}, \boldsymbol{\phi}) = -\ln p_{\boldsymbol{\theta}}(\mathbf{x}_i) + D_{KL} \left( q_{\boldsymbol{\phi}}(\mathbf{z}|\mathbf{x}_i) || p_{\boldsymbol{\theta}}(\mathbf{z}|\mathbf{x}_i) \right)$$

maximize likelihood of observed data (minimize reconstruction error)

and

minimize difference of approximation  $q_{\phi}(\mathbf{z}|\mathbf{x}_i)$  to exact posterior  $p_{\theta}(\mathbf{z}|\mathbf{x}_i)$ 

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can be interpreted as regularizer

corresponds to maximizing evidence lower bound (ELBO), i.e., maximizing lower bound of probability to generate real data sample:

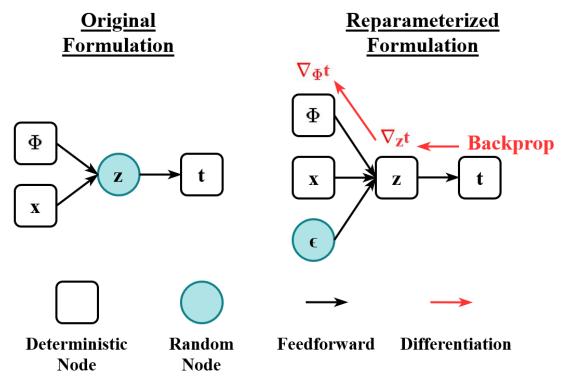
$$\ln p_{\boldsymbol{\theta}}(\boldsymbol{x}_i) \ge \ln p_{\boldsymbol{\theta}}(\boldsymbol{x}_i) - D_{KL}\left(q_{\boldsymbol{\phi}}(\boldsymbol{z}|\boldsymbol{x}_i)||p_{\boldsymbol{\theta}}(\boldsymbol{z}|\boldsymbol{x}_i)\right) = E_{\boldsymbol{z} \sim q_{\boldsymbol{\phi}}(\boldsymbol{z}|\boldsymbol{x}_i)}\left[\ln \frac{p_{\boldsymbol{\theta}}(\boldsymbol{x}_i, \boldsymbol{z})}{q_{\boldsymbol{\phi}}(\boldsymbol{z}|\boldsymbol{x}_i)}\right]$$
non-negative

## Reparameterization Trick

ightarrow gradient descent according to  $oldsymbol{ heta}$  and  $oldsymbol{\phi}$ 

issue: not readily possible for  $\phi$  (expecatation over z, which is sampled from  $q_{\phi}$ )

 $\rightarrow$  reparametrization to the rescue: express randomness in z by independent auxiliary variable  $\varepsilon$ 



from wikipedia

## Gaussian Approximation

Posterior and approximation

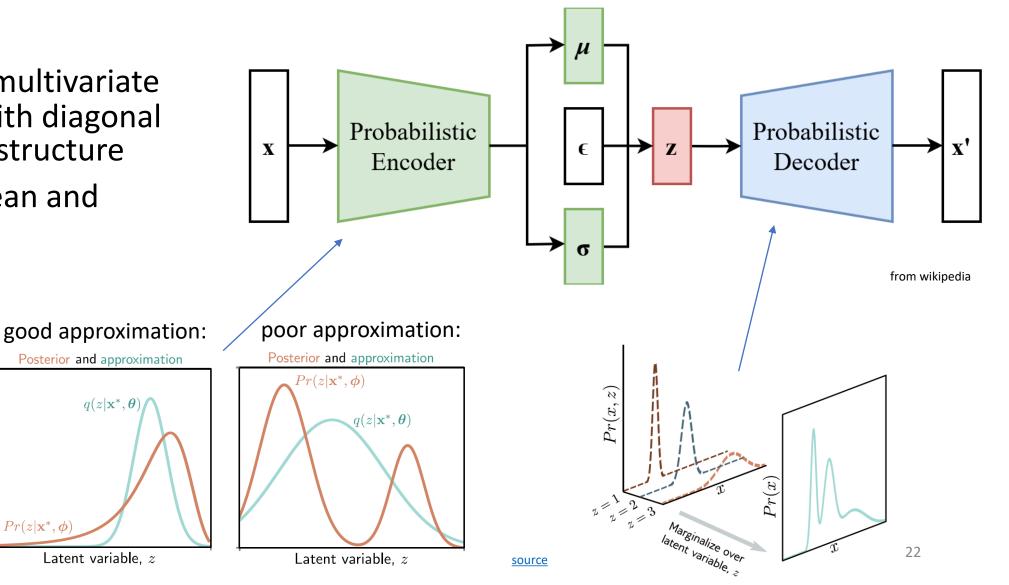
Latent variable, z

 $Pr(z|\mathbf{x}^*, \boldsymbol{\phi})$ 

 $q(z|\mathbf{x}^*, \boldsymbol{\theta})$ 

e.g.,  $q_{\phi}$  as multivariate Gaussian with diagonal covariance structure

→ learn mean and variance



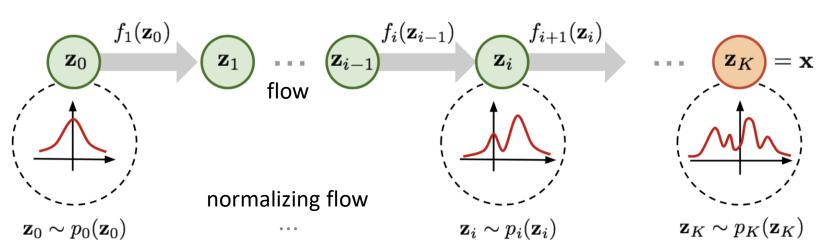
## Flow-Based Methods

## Normalizing Flows

idea: mapping of a simple probability distribution (often, standard normal distribution) into a complex one by sequence of invertible transformations (repeatedly applying the change-of-variable technique)

$$f(\mathbf{z}') = f(\mathbf{z}) \left| \det \frac{\delta f^{-1}}{\delta \mathbf{z}'} \right| = f(\mathbf{z}) \left| \det \frac{\delta f}{\delta \mathbf{z}} \right|^{-1}$$
$$\ln p_K(\mathbf{z}_K) = \ln p_0(\mathbf{z}_0) - \sum_{k=1}^K \ln \left| \det \frac{\delta f_k}{\delta \mathbf{z}_{k-1}} \right|$$

log-likelihood:



## Usage in Generative Models

training: estimate maximum likelihood of normalizing flow (log-likelihood of last slide) by gradient descent (learn parameters  $\theta$  of transformations  $f_{\theta}^{-1}$ , e.g., to let  $p_0(\mathbf{z})$  be Gaussian)

inference: sample from simple distribution  $p_0(\mathbf{z})$  and transform it back to data distribution  $p_K(\mathbf{x})$  via  $f_\theta$ 

#### advantages:

- instead of simple distributions like Gaussians, allow more complex latent encodings: real-world distributions usually much more complicated
- exact likelihood estimation (VAEs and diffusion models only return lower bound): allows density estimation (e.g., for anomaly detection)

#### Invertible Neural Networks

neural networks representing invertible/bijective functions can be used for normalizing flow transformations

- forward transformation to generate samples
- backward transformation to evaluate likelihoods

need for specialized architectures to construct reversible transform (e.g., affine coupling layers)

## Diffusion Models

## Idea

training: distort training data by successively adding random noise, then learn to reverse this process (denoising)

generation: sample random noise and run through the learned denoising process

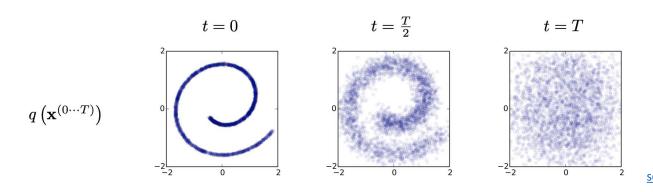
advantages: easy to train, produce high-quality/realistic samples can be interpreted as special case of hierarchical VAE (one latent variable generates another) with fixed encoder and latent space of same size as the data  $\rightarrow$  more sophisticated latent space than just Gaussian mixture in VAE

#### Forward Process

Markov chain of diffusion steps to slowly add Gaussian noise to data (inspired by non-equilibrium thermodynamics):

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

- with variance schedule  $\beta_1, ..., \beta_T$  (hyperparameters, increasing with t)
- large T and small  $\beta_t \rightarrow$  same functional form for forward and reverse processes, ending up with isotropic Gaussian distribution for  $x_T$



## Reparametrization

conditional Gaussian distributions at each t:

sample 
$$\epsilon \sim \mathcal{N}(0, \mathbf{I})$$
 and set  $\mathbf{x}_t = \sqrt{1 - \beta_t} \mathbf{x}_{t-1} + \sqrt{\beta_t} \epsilon$ 

nice property: possible to directly sample  $x_t$  conditioned on  $x_0$  (no need to apply q repeatedly)

$$\begin{aligned} \boldsymbol{x}_t &= \sqrt{\bar{\alpha}_t} \boldsymbol{x}_0 + \sqrt{1 - \bar{\alpha}_t} \boldsymbol{\epsilon} \\ q(\boldsymbol{x}_t | \boldsymbol{x}_0) &= \mathcal{N} \big( \boldsymbol{x}_t; \, \sqrt{\bar{\alpha}_t} \boldsymbol{x}_0, (1 - \bar{\alpha}_t) \mathbf{I} \big) \end{aligned}$$
 with  $\alpha_t = 1 - \beta_t$ ,  $\bar{\alpha}_t = \prod_{s=1}^t \alpha_s$ 

conditioning on  $x_0$  also allows to handle  $q(x_{t-1}|x_t,x_0)$ 

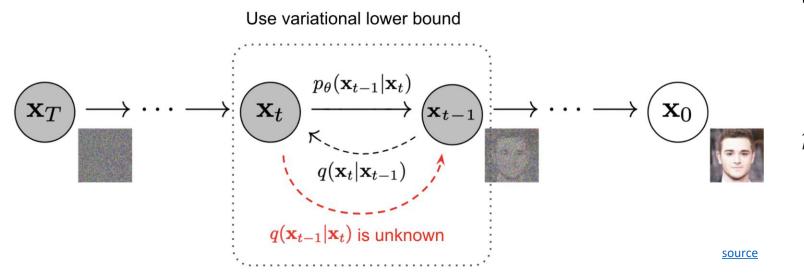
#### Reverse Process

to generate new data samples, one needs to learn to reverse the diffusion process (starting from pure noise): neural network learning to gradually denoise data

overall loss as sum of losses for each time step t

for each t:  $D_{KL}$  between two Gaussians (closed form)  $q(x_{t-1}|x_t,x_0)$  and  $p_{\theta}(x_{t-1}|x_t)$ 

→ corresponds to VAE loss: maximizing ELBO



time-dependent Gaussian parameters:

$$p(\mathbf{x}_T) = \mathcal{N}(\mathbf{x}_T; 0, \mathbf{I})$$

$$p_{\theta}(\mathbf{x}_{0:T}) = p(\mathbf{x}_T) \prod_{t=1}^{T} p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t)$$

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

#### Noise Prediction

reparametrization allows to learn added noise instead of Gaussian parameters:

Algorithm 1 Training	Algorithm 2 Sampling
1: repeat 2: $\mathbf{x}_0 \sim q(\mathbf{x}_0)$ 3: $t \sim \mathrm{Uniform}(\{1,\ldots,T\})$ 4: $\boldsymbol{\epsilon} \sim \mathcal{N}(0,\mathbf{I})$ 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	1: $\mathbf{x}_{T} \sim \mathcal{N}(0, \mathbf{I})$ 2: <b>for</b> $t = T, \dots, 1$ <b>do</b> 3: $\mathbf{z} \sim \mathcal{N}(0, \mathbf{I})$ if $t > 1$ , else $\mathbf{z} = 0$ 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: <b>end for</b> 6: <b>return</b> $\mathbf{x}_{0}$

L2-loss (MSE) between true and predicted Gaussian noise at time step t use position embeddings (as network parameters are shared across time)

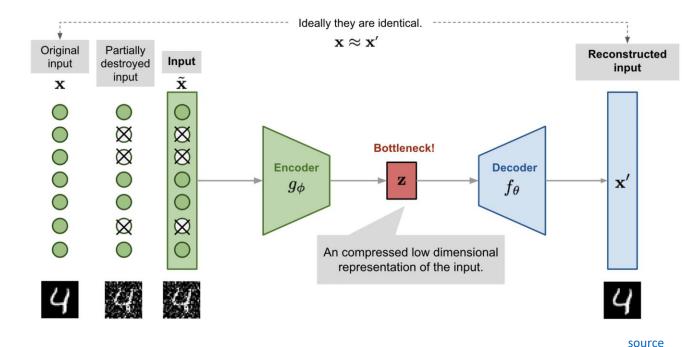
diffusion models can be interpreted as chain of denoising autoencoders (also connected to score-based generative modeling via Langevin dynamics)

## Denoising Autoencoder

goal: avoid overfitting and improve robustness of plain autoencoder

learn to remove noise of distorted input  $\tilde{x} \rightarrow$  restore original input x

similar to dropout



differences of diffusion models to typical denoising autoencoders:

- no bottleneck (care about output here, not internal representation): latent space with high dimensionality (same as original data)
- handle many different noise levels with single set of shared parameters

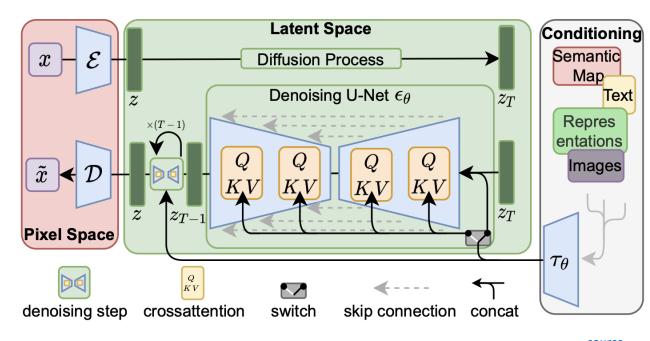
## Latent Diffusion Model

add noise to latent representation rather than raw data

→ significant speedup

diffusion models highly flexible in terms of architecture: only require same input and output dimensionality (autoencoder-like)

- often, (convolutional) U-Net architectures (skip connections)
- but also (vision) transformers possible (e.g., <u>DiT</u>)



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use of attention mechanism for flexible conditioning

## Conditioned Generation

#### Conditional GANs

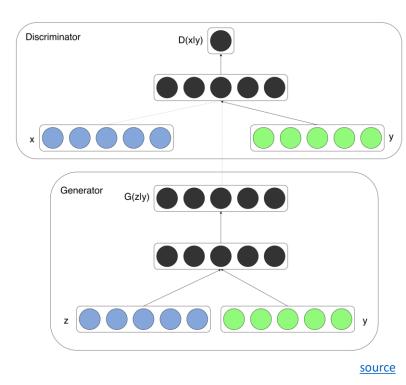
as discussed so far, generative methods give no control over what kind of data is generated (limited usability)

→ need for conditional approach (e.g., conditioning on describing text)

#### example GANs:

transform usual GAN to conditional model by feeding extra information y (e.g., class labels) as additional input layer into both generator and discriminator

$$L(\mathbf{x}_{i}) = E_{\mathbf{x} \sim p_{r}(\mathbf{x})} [\ln D(\mathbf{x}_{i}|y_{i})] + E_{\mathbf{x} \sim p_{g}(\mathbf{x})} [\ln (1 - D(\mathbf{x}_{i}|y_{i}))]$$



## **Guided Diffusion**

ways to condition on class information in diffusion process:

- classifier guidance: perturbation of classconditional diffusion model by separately trained classifier model  $p_{\theta}(y|x_t)$ 
  - $\widehat{\mu}_{\theta}(x_t|y) = \mu_{\theta}(x_t|y) + s \cdot \Sigma_{\theta}(x_t|y) \cdot \nabla_{x_t} \log p_{\theta}(y|x_t)$
  - guidance can also be free-form text, e.g., from <a href="CLIP">CLIP</a> model
- classifier-free guidance: randomly replace label in class-conditional diffusion model with null label during training

extrapolate in direction of conditioned model during sampling:  $\hat{\epsilon}_{\theta}(x_t|y) \neq \hat{\epsilon}_{\theta}(x_t|\emptyset) + s \cdot (\hat{\epsilon}_{\theta}(x_t|y) - \hat{\epsilon}_{\theta}(x_t|\emptyset))$ 

similar idea as softmax temperature in auto-regressive LLMs

tradeoff between diversity (unconditioned) and fidelity (guidance)



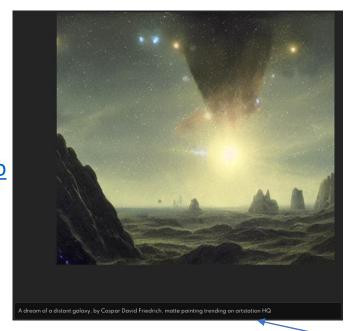
"Pembroke Welsh corgi"

## Text-to-Image

plenty of applications: <u>DALL-E 2</u>, <u>Stable Diffusion</u>, <u>ImageGen</u>, ...

also text-to-speech (e.g., <u>VALL-E</u>), text-to-video (e.g., <u>Make-A-Video</u>), ...

web app for Stable Diffusion: DreamStudio



inpainting example (GLIDE):



"zebras roaming in the field"

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prompt

#### Literature

#### papers:

- variational autoencoder
- normalizing flows
- GAN
- denoising diffusion, latent diffusion



## Movie-like Intelligence

**emergent capabilities** of complex systems difficult to foresee

mini examples in contemporary ML:

- large language models
- multi-agent reinforcement learning

philosophical: emotions and consciousness in humans may also have occurred as emergent capabilities (But that does not mean the same will happen with AI.)

ideas for paths toward general intelligence:

- reward is enough
  - small-world/scale-free networks