



# CS 839: Foundation Models

## **Diffusion Models**

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# Announcements

- **Logistics:**

- Project info out
- Come chat about presentation/project proposals.
- Class roadmap:

Thursday Oct. 30	Diffusion Models
Tuesday Nov. 4	Scaling & Scaling Laws
Thursday Nov. 6	Security, Privacy, Toxicity + Future Areas

# Outline

- **Generative Models Overview**

- Basic idea, complexity challenges, overview of major image generation techniques, intuitions

- **Normalizing Flows & GANs**

- Normalizing flow transformations, training, sampling, GAN generators, discriminators, training

- **Diffusion Models**

- Overall intuition, score-based training, controlling and latent space formulations, extensions

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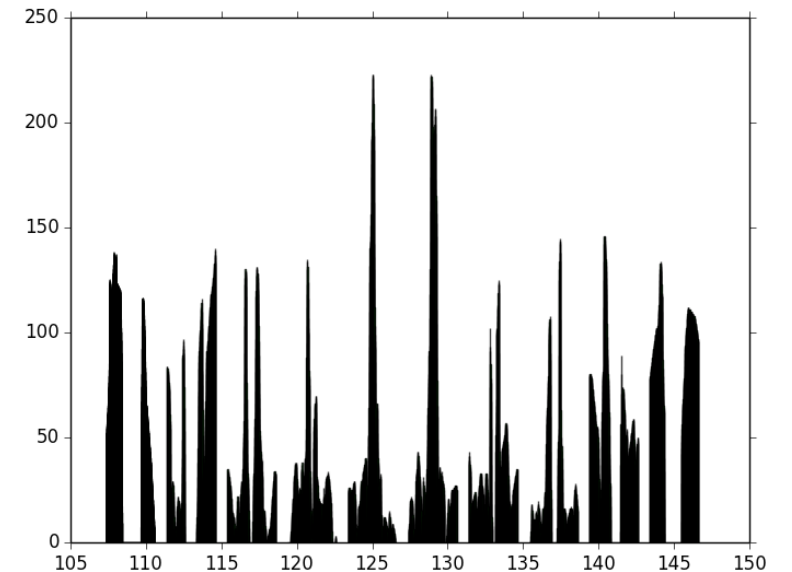
- Overall intuition, score-based training, controlling and latent space formulations, extensions

# Goal: Learn a Distribution

- Want to estimate  $p_{\text{data}}$  from samples

$$x^{(1)}, x^{(2)}, \dots, x^{(n)} \sim p_{\text{data}}(x)$$

- Useful abilities to have:
  - **Inference**: compute  $p(x)$  for some  $x$
  - **Sampling**: obtain a sample from  $p(x)$
- As always need efficiency for this too...



# Directly Modeling the Distribution

- Want to estimate  $p_{\text{data}}$  from samples

$$x^{(1)}, x^{(2)}, \dots, x^{(n)} \sim p_{\text{data}}(x)$$

- One straightforward idea: **parametrize the pdf of the distribution**. To train, maximize the log likelihood

$$\max_{\theta} \sum_{i=1}^N \log p_{\theta}(x_i).$$

- However, we'll face some challenges...
  - Why? Both training and inference can be complex

# Goal: Learn a Distribution

- Want to estimate  $p_{\text{data}}$  from samples

$$x^{(1)}, x^{(2)}, \dots, x^{(n)} \sim p_{\text{data}}(x)$$

- Let's set

$$p_{\theta}(x) = \frac{1}{Z} \exp(f_{\theta}(x))$$

Energy function



- Have to deal with the normalizing **partition function Z**,

$$Z_{\theta} = \int \exp(f_{\theta}(x)) dx$$

**Usually intractable!**

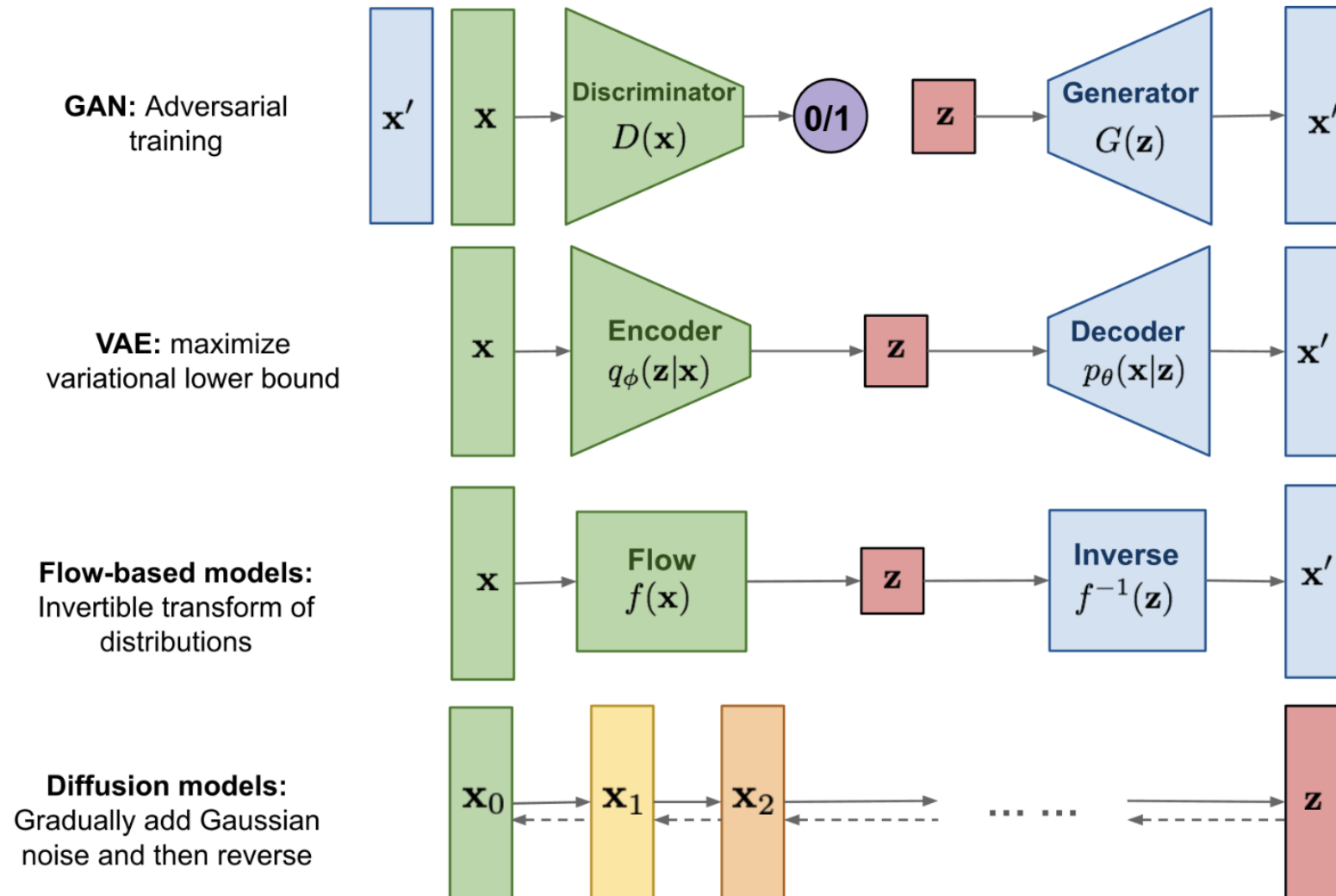
# Getting Around the Partition Function

All gen. modeling techniques must deal with this. How?

- Avoid modeling the pdf explicitly
  - → **GANs**
- Choose special choices of p/f that keeps  $Z$  tractable
  - → Certain **normalizing flows**
- Use approximations
  - → **VAEs**, using ELBO-style bounds
- Obtain training objectives that sidestep maximum likelihood
  - → **GANs, score-based diffusion models**



# Generative Modeling Approaches

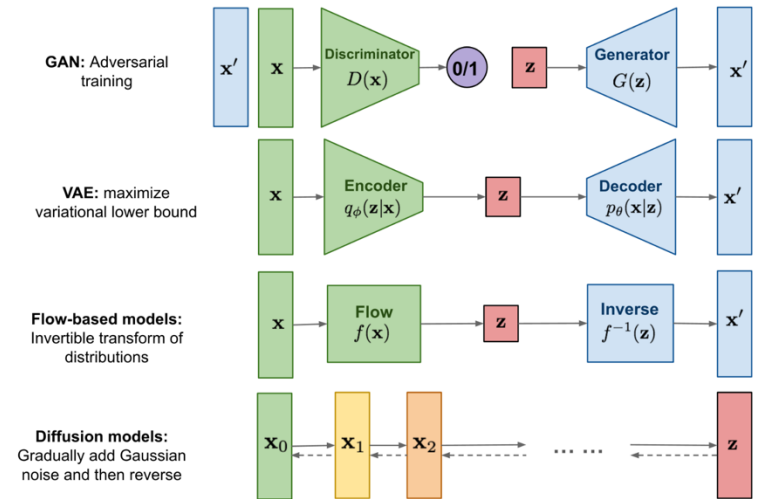


# Generative Modeling Intuitions

We can think of GMs as doing two things:

- “Mapping” a **simple** (fake) distribution into a **complex** (real) distribution
  - Why? Sample from simple distribution, then transform with learned map
  - “**Latent space**” interpretation
- Learning to undo noise or undo a particular transformation
  - Related to self-supervised learning

Combine with previous training considerations to get various techniques



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**Break & Questions**

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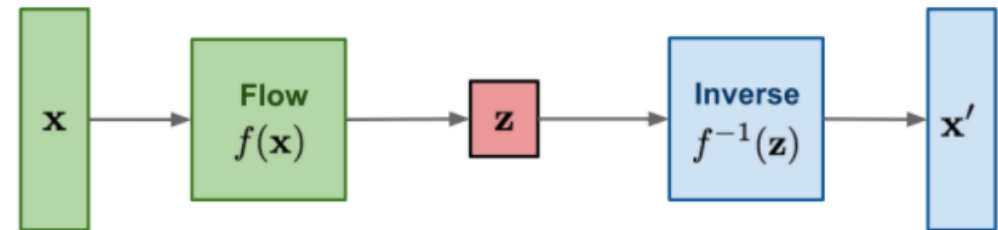
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# Flow Models

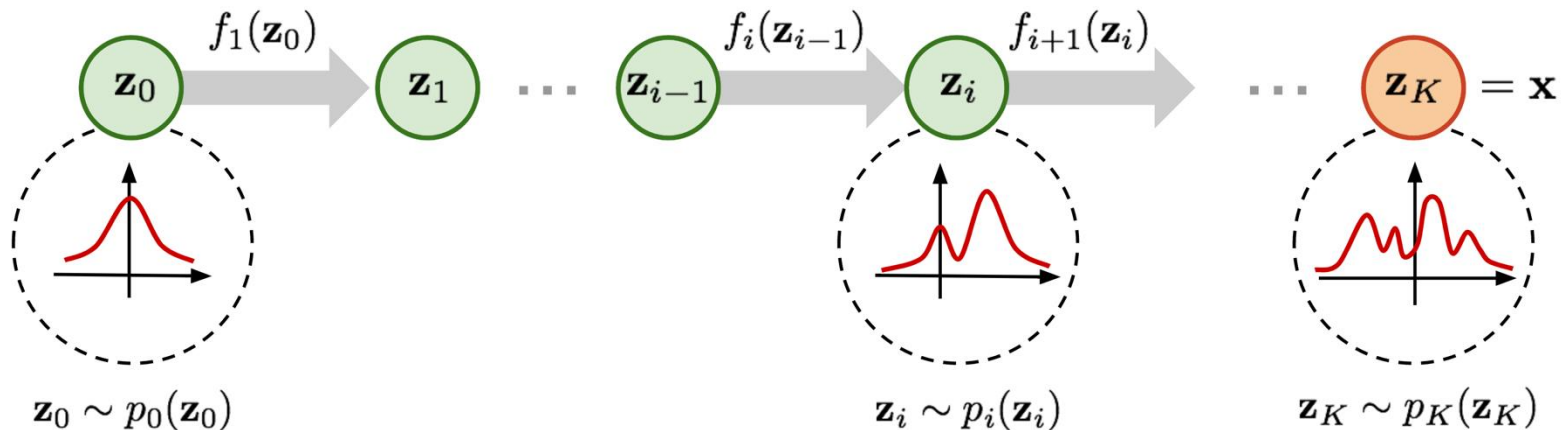
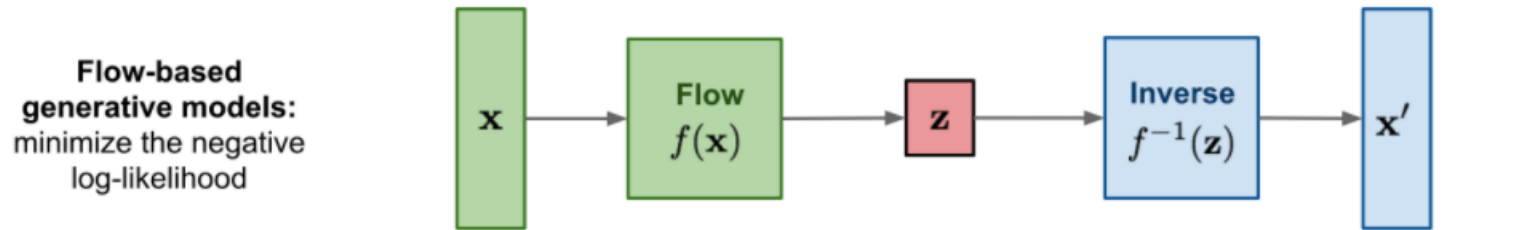
- Want to fit  $p_{\theta}(x)$ , as we described
- Some goals:
  - Good fit for the data
  - Computing a probability: the actual value of  $p_{\theta}(x)$  for some  $x$
  - Ability to sample
  - Also: a **latent representation**
- Won't model  $p_{\theta}(x)$  directly... instead we'll get some latent variable  $z$

Flow-based  
generative models:  
minimize the negative  
log-likelihood



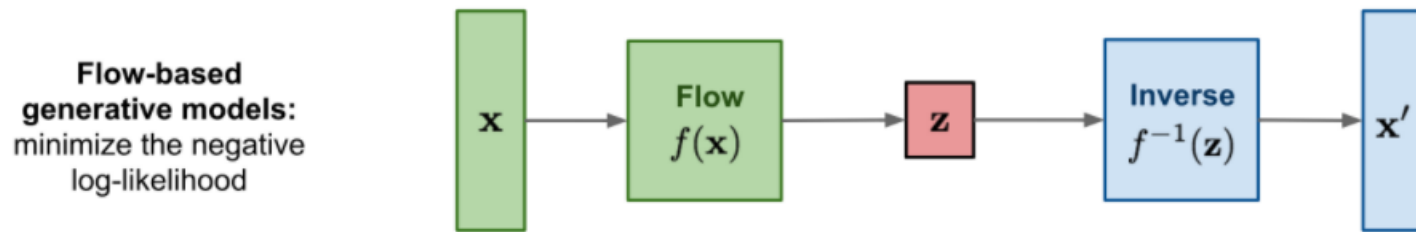
# Flow Models

- **Key idea:** transform a simple distribution to complex
  - Use a chain of transformations (the “flow”)



# Flow Models

- **Key idea:** transform a simple distribution to complex
  - Use a chain of invertible transformations (the “flow”)



- How to sample?
  - Sample from  $Z$  (the latent variable)---has a simple distribution that lets us do it: Gaussian, uniform, etc.
  - Then run the sample  $z$  through the inverse flow to get a sample  $x$
- How to train? Let's see...



# Flow Models: Density Relationships

- **Key idea:** transform a simple distribution to complex
  - Use a chain of transformations (the “flow”)
- How does each transformation affect the density  $p$ ?

Latent variable      Transformation

$$z = f_{\theta}(x)$$
$$p_{\theta}(x) dx = p(z) dz$$
$$p_{\theta}(x) = p(f_{\theta}(x)) \left| \frac{\partial f_{\theta}(x)}{\partial x} \right|$$

Determinant of Jacobian matrix



# Flow Models: Training

- **Key idea:** transform a simple distribution to complex
  - Use a chain of transformations (the “flow”)
- How does training change?
  - **Idea:** might be easier to optimize  $p_Z$

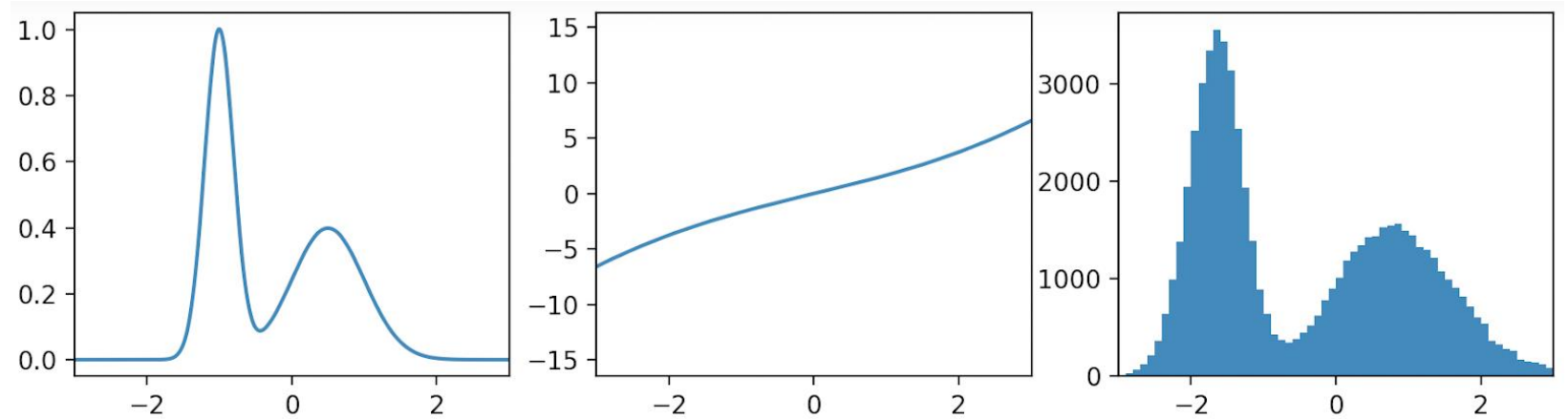
$$\max_{\theta} \underbrace{\sum_i \log p_{\theta}(x^{(i)})}_{\text{Maximum Likelihood}} = \max_{\theta} \sum_i \log p_Z(\underbrace{f_{\theta}(x^{(i)})}_{\text{Latent variable version}}) + \log \underbrace{\left| \frac{\partial f_{\theta}}{\partial x}(x^{(i)}) \right|}_{\text{Determinant of Jacobian matrix}}$$

Can extend to many chained transformations...

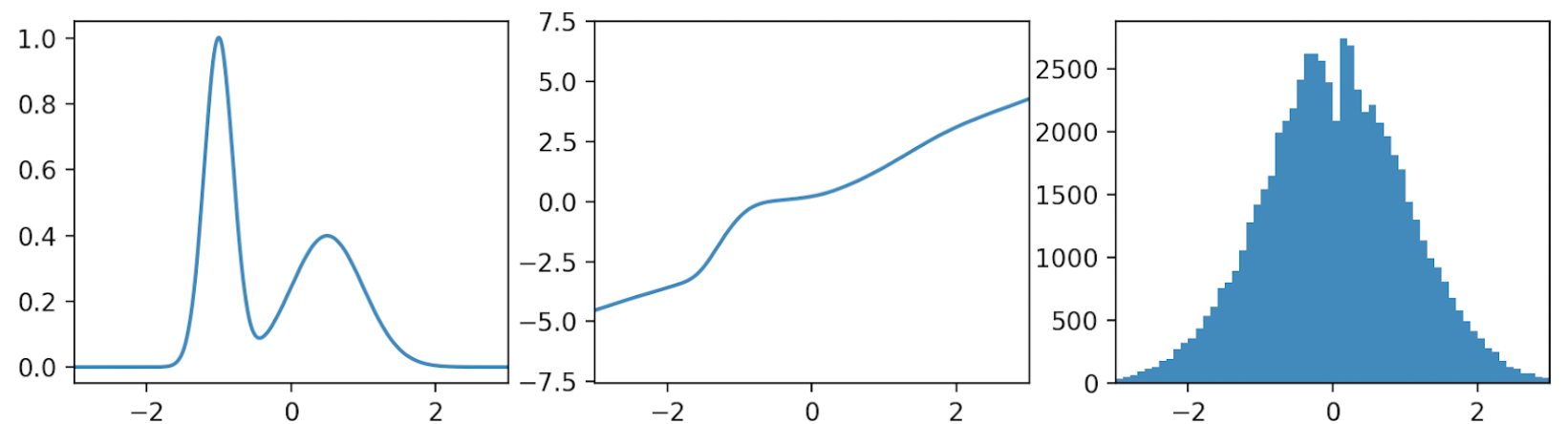
# Flows: Example

- Flow to a Gaussian (right)

- Before training:

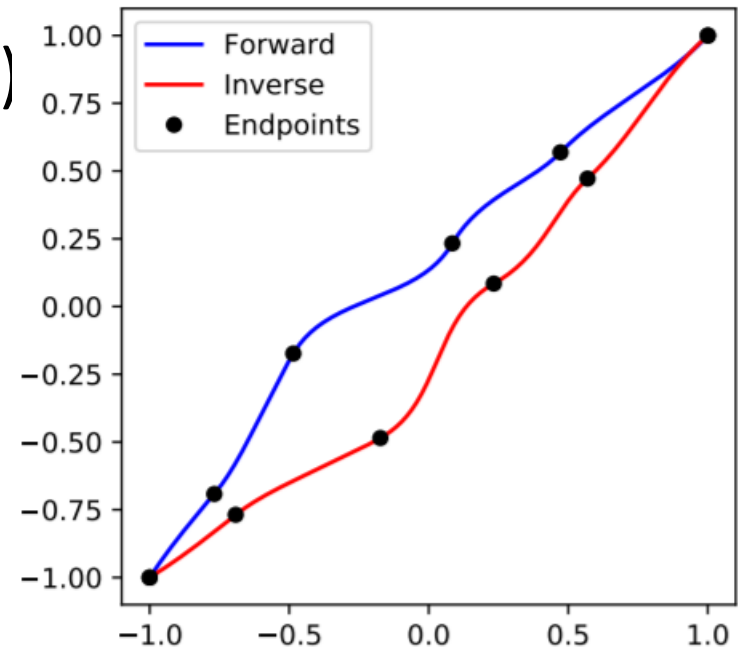


- After training:



# Flows: Transformations

- What kind of  $f$  transformations should we use?
- Many choices:
  - Affine:  $f(x) = A^{-1}(x - b)$
  - Elementwise:  $f(x_1, \dots, x_d) = (f(x_1), \dots, f(x_d))$
  - Splines:
- Desirable properties:
  - Invertible
  - Differentiable (forward and inverse)



(a) Forward and inverse transformer

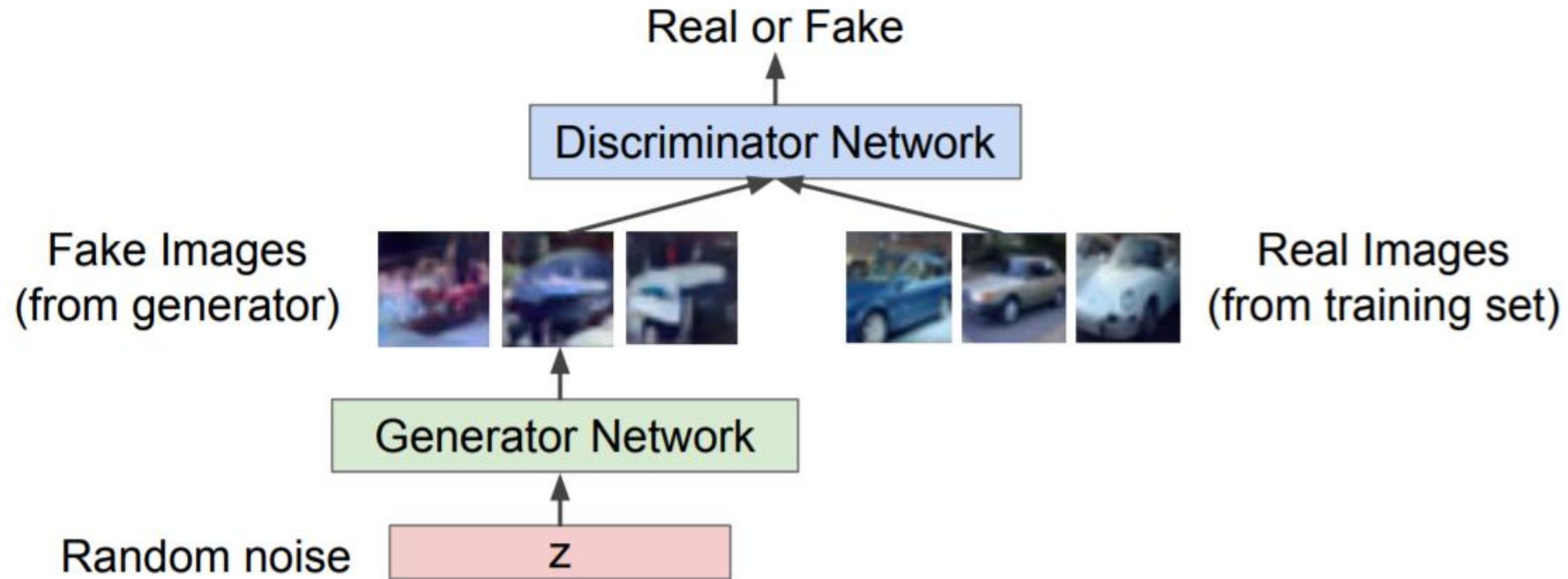
# GANs: Generative Adversarial Networks

- So far, we've been modeling the density...
  - What if we just want to get high-quality samples?
- GANs do this. Based on a clever idea:
  - Art forgery: very common through history
  - Left: original
  - Right: forged version
  - Two-player game. **Forger** wants to pass off the forgery as an original; **investigator** wants to distinguish forgery from original



# GANs: Basic Setup


- Let's set up networks that implement this idea:
  - Discriminator network: like the **investigator**
  - Generator network: like the **forgery**




# GAN Training: Discriminator

- How to train these networks? Two sets of parameters to learn:  $\theta_d$  (**discriminator**) and  $\theta_g$  (**generator**)
- Let's fix the generator. What should the discriminator do?
  - Distinguish fake and real data: binary classification.
  - Use the cross entropy loss, we get

$$\max_{\theta_d} \mathbb{E}_{x \sim p_{\text{data}}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$



**Real data, want  
to classify 1**



**Fake data, want  
to classify 0**

# GAN Training: Generator & Discriminator

- How to train these networks? Two sets of parameters to learn:  $\theta_d$  (**discriminator**) and  $\theta_g$  (**generator**)
- This makes the discriminator better, but also want to make the generator more capable of fooling it:
  - Minimax game! Train jointly.

$$\min_{\theta_g} \max_{\theta_d} \mathbb{E}_{x \sim p_{\text{data}}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$

↑  
**Real data, want  
to classify 1**

↑  
**Fake data, want  
to classify 0**

# GAN Training: Alternating Training

- So we have an optimization goal:

$$\min_{\theta_g} \max_{\theta_d} \mathbb{E}_{x \sim p_{\text{data}}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$

- Alternate training:

- **Gradient ascent**: fix generator, make the discriminator better:

$$\max_{\theta_d} \mathbb{E}_{x \sim p_{\text{data}}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$

- **Gradient descent**: fix discriminator, make the generator better

$$\min_{\theta_g} \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$

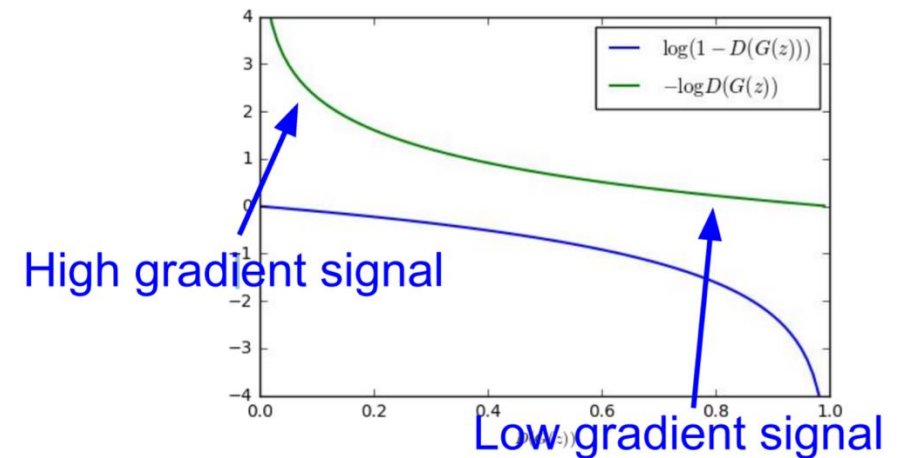


# GAN Training: Issues

- Training often not stable
- Many tricks to help with this:
  - Replace the generator training with

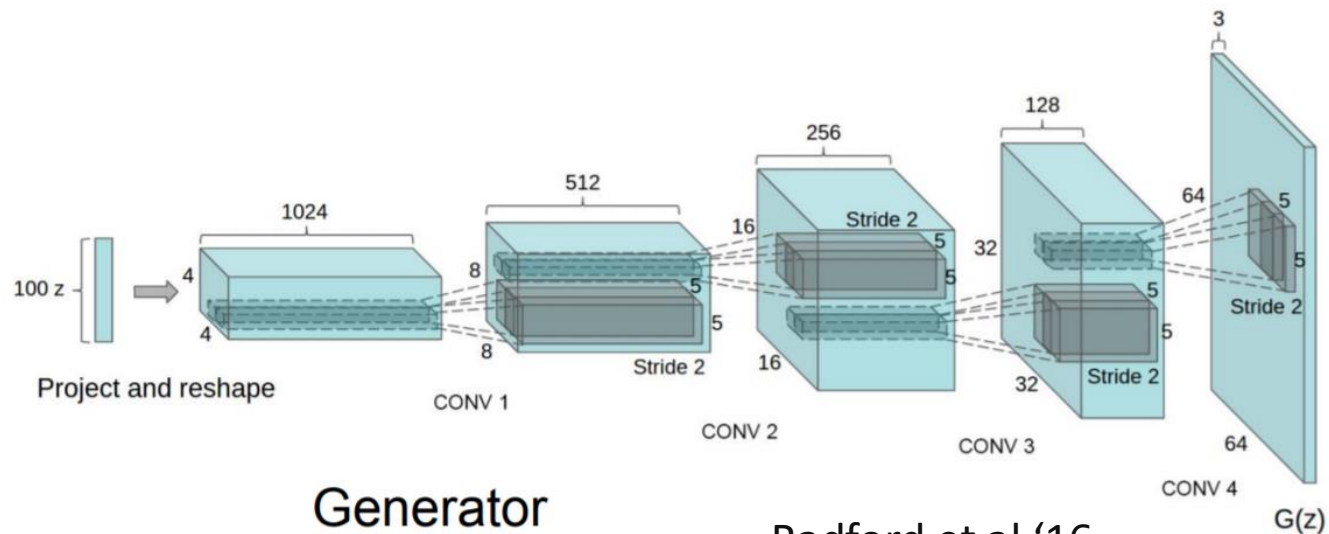
$$\max_{\theta_g} \mathbb{E}_{z \sim p(z)} \log(D_{\theta_d}(G_{\theta_g}(z)))$$

- Better gradient shape
  - Choose number of alt. steps carefully
- Can still be challenging.



# GAN Architectures

- So far we haven't commented on what the networks are
- **Discriminator**: image classification, use a **CNN**
- What should **generator** look like
  - Input: noise vector  $z$ . Output: an image (ie, volume 3 x width x height)
  - Can just reverse our CNN pattern...



# GANs: Example

- From Radford's paper, with 5 epochs of training:







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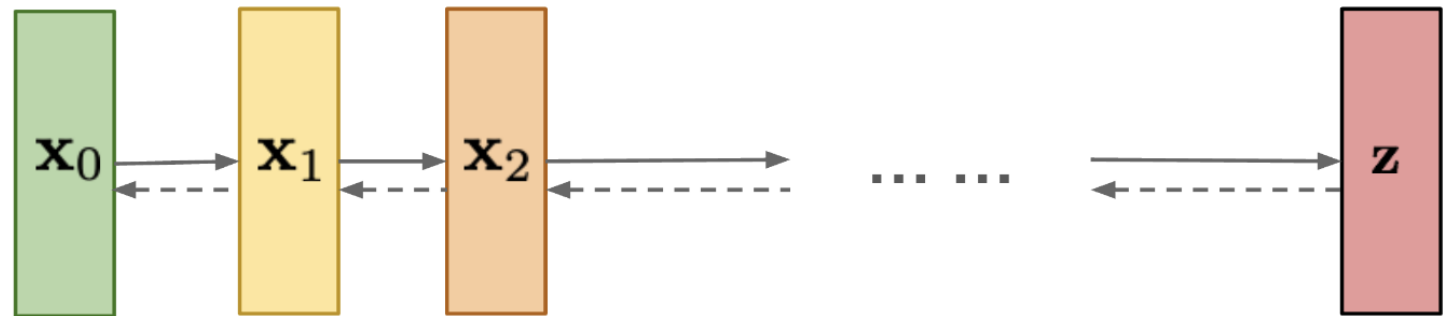
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# Diffusion Models Idea

- Let's return to something that looks like a normalizing flow,

**Diffusion models:**  
Gradually add Gaussian  
noise and then reverse



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- Really a large family of techniques that share some common properties
  - But have been derived from different starting principles / desired properties

# Score-Based Generative Models

- How do we avoid running into the partition function?
- Let's not model the pdf
- Instead, model the “**score**”

$$\nabla_{\mathbf{x}} \log p(\mathbf{x})$$

- Score: gradient of the log likelihood with respect to the data.
- Goal: train  $s$  such that

$$\mathbf{s}_{\theta}(\mathbf{x}) = \nabla_{\mathbf{x}} \log p_{\theta}(\mathbf{x}) :$$

# Score-Based Generative Models

Instead, model the “score”


$$\nabla_{\mathbf{x}} \log p(\mathbf{x})$$

Goal: train  $s$  such that

$$\mathbf{s}_{\theta}(\mathbf{x}) = \nabla_{\mathbf{x}} \log p_{\theta}(\mathbf{x}) :$$

- Why does this avoid the partition function?
- Let's plug in our energy-based function from earlier. We get:

Gradient w.r.t. **x**, not  $\theta$


$$\mathbf{s}_{\theta}(\mathbf{x}) = \nabla_{\mathbf{x}} \log p_{\theta}(\mathbf{x}) = -\nabla_{\mathbf{x}} f_{\theta}(\mathbf{x}) - \underbrace{\nabla_{\mathbf{x}} \log Z_{\theta}}_{=0} = -\nabla_{\mathbf{x}} f_{\theta}(\mathbf{x}).$$



# Training & Inference for Score-Based Models

- Training: can directly run M.S.E. as a loss,

$$\mathbb{E}_{p(\mathbf{x})} [\|\nabla_{\mathbf{x}} \log p(\mathbf{x}) - \mathbf{s}_{\theta}(\mathbf{x})\|_2^2]$$

- We usually can't access the left hand term, but techniques for training despite this
- Inference: special methods that can sample, like Langevin dynamics

$$\mathbf{x}_{i+1} \leftarrow \mathbf{x}_i + \epsilon \nabla_{\mathbf{x}} \log p(\mathbf{x}) + \sqrt{2\epsilon} \mathbf{z}_i$$



Sample  
Iterates



Learned  
score function

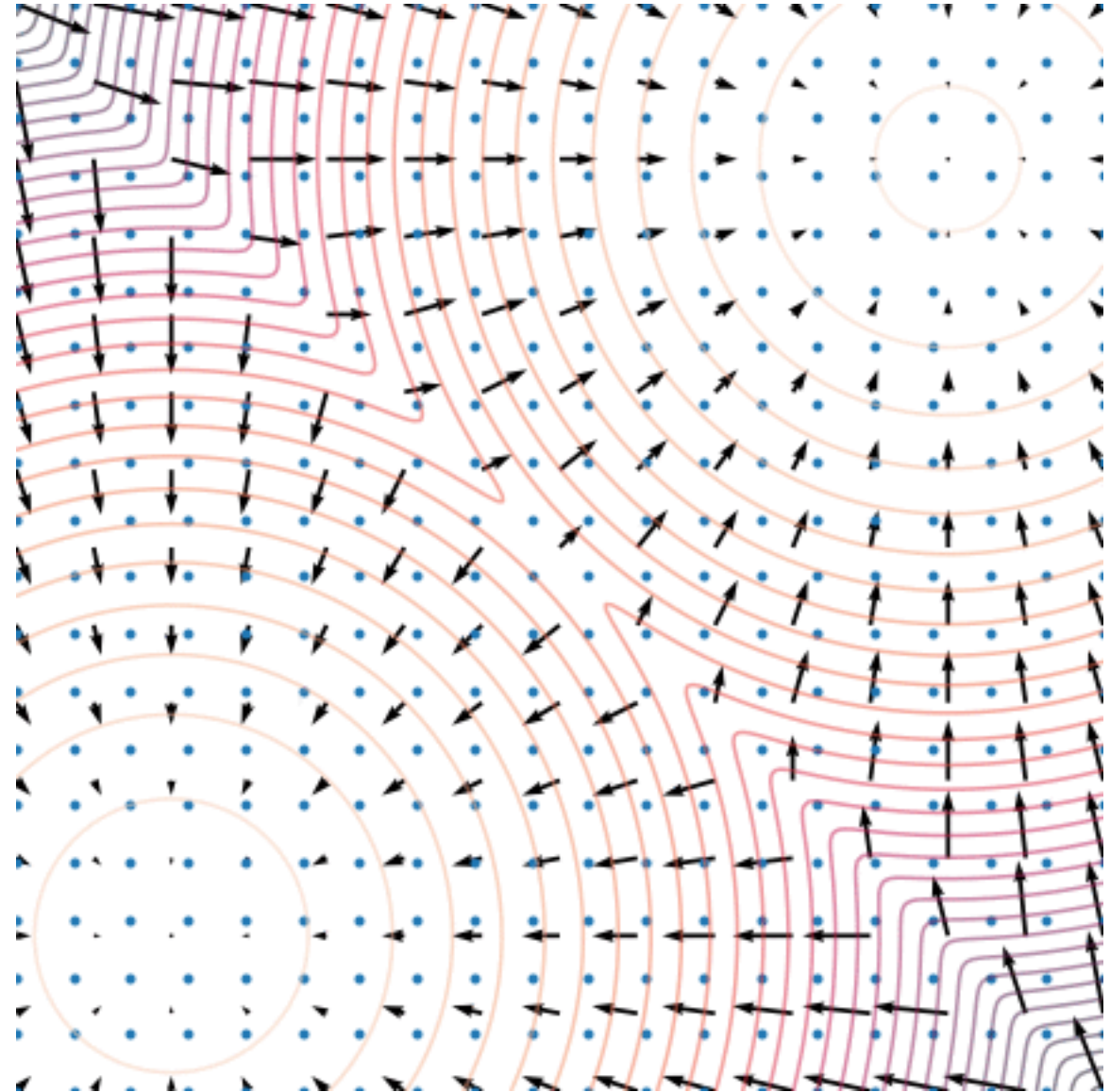


Noise

# Training & Inference for Score-Based Models

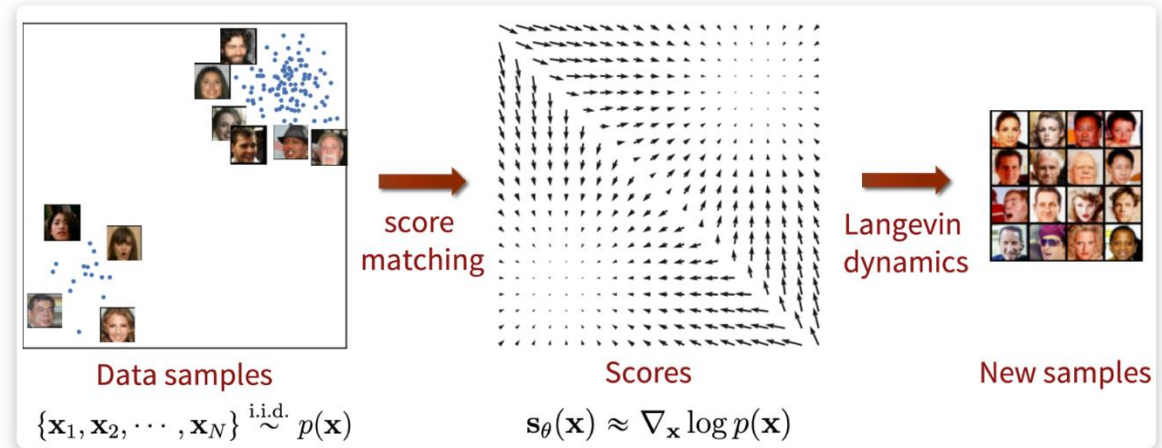
- **Visual example**

- Distribution: mixture of two Gaussians
- Arrows: given by our score function, point to high density regions
- Source: <https://yang-song.net/blog/2021/score/>

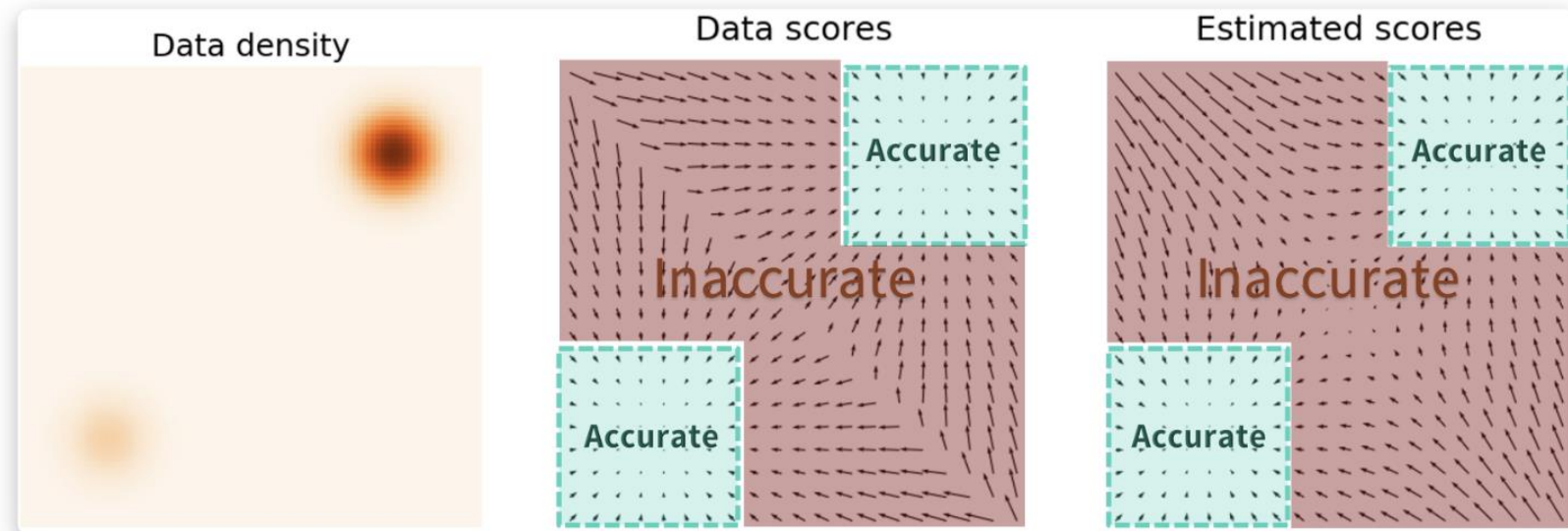


# Score-Based → Denoising Diffusion Models

- Our story so far is

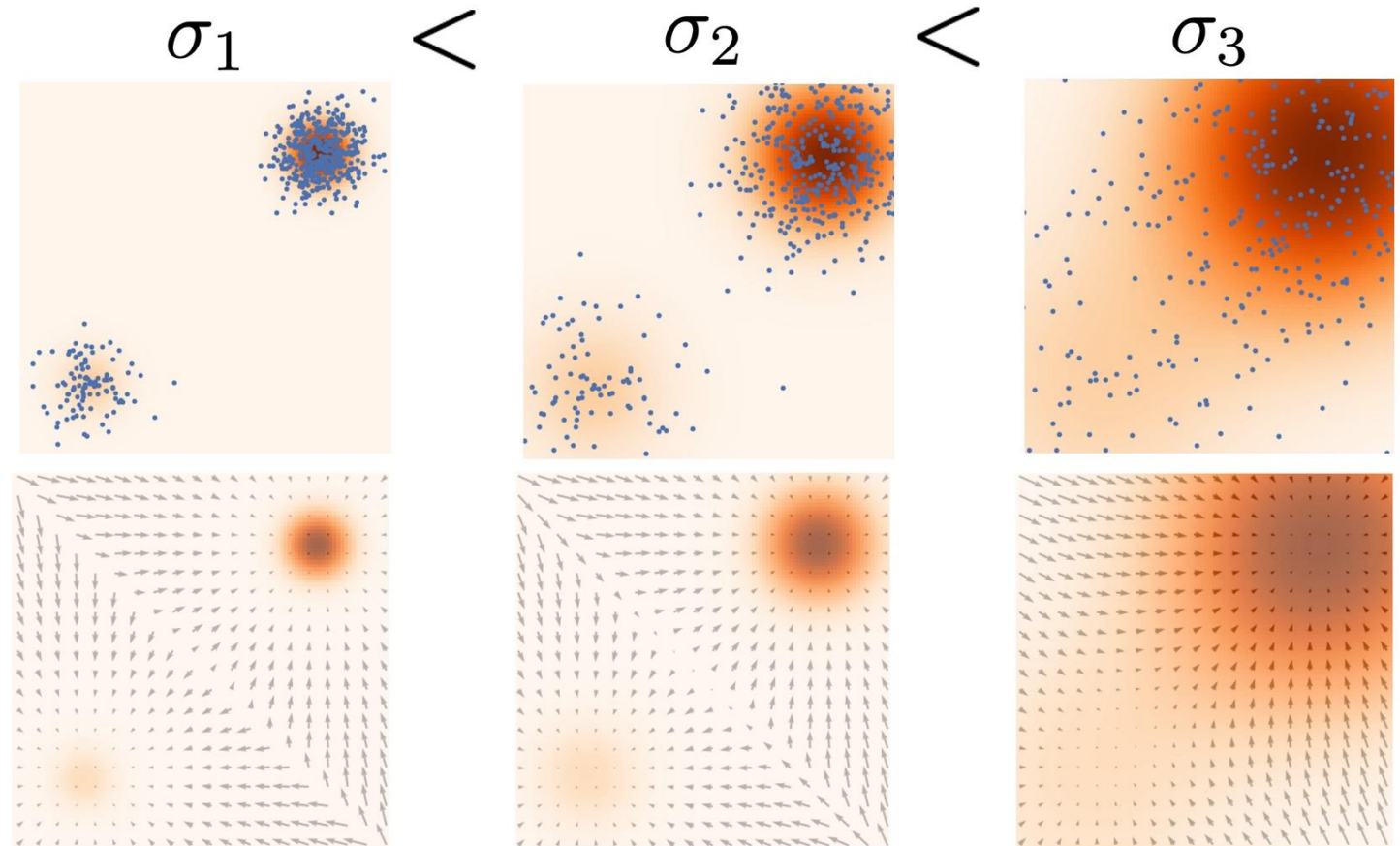


- But, this leads to inaccurate modeling in low-prob regions:



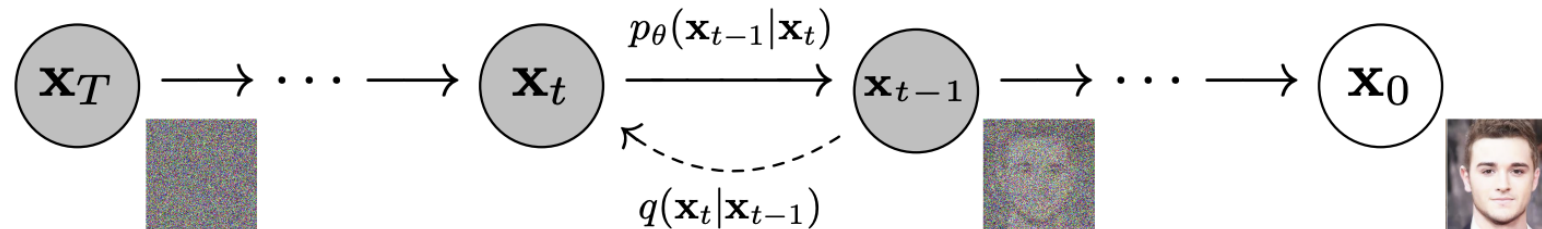
# Score-Based → Denoising Diffusion Models

- Solution: perturb the density with noise
  - To ensure accurate modeling in more regions
  - In particular, noise at multiple scales



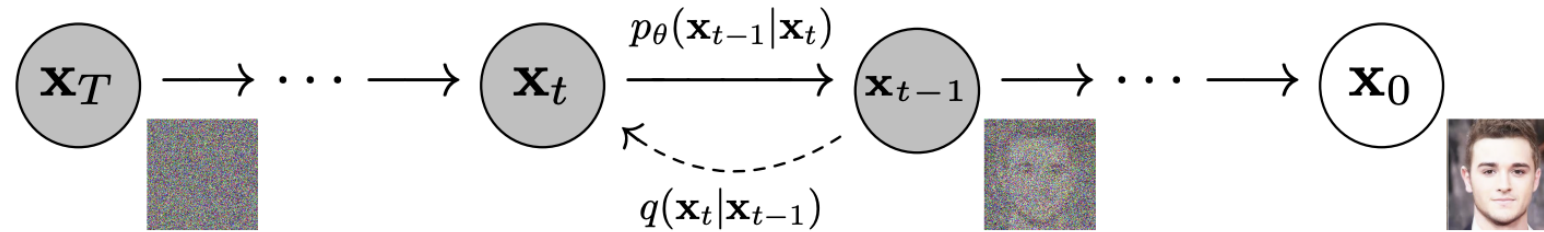
# Score-Based → Denoising Diffusion Models

- So far, “noise” showed up in a few places, but not in a strictly connected way
  - Train model with score matching
  - Sample with Langevin dynamics (which includes noise)
  - Use noise perturbation to train better
- Denoising diffusion models **directly** use noise in both training and inference



# Diffusion Models

- Basic graphical model



Ho et al '20

- Can easily set up the noising process,

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

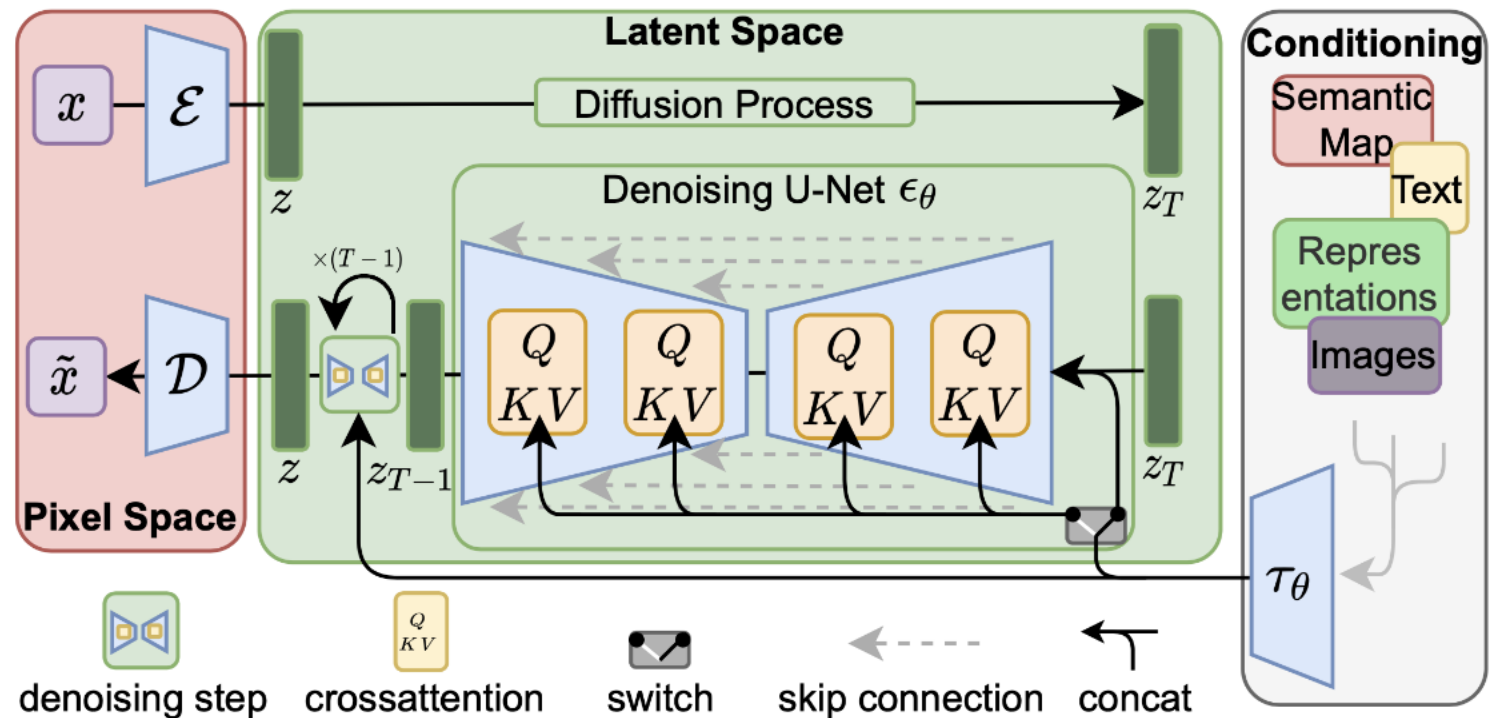
- To sample, directly compute from reverse, i.e.,  $q(\mathbf{x}_{t-1}|\mathbf{x}_t)$ 
  - Simple, nice parametrizations in Ho et al '20.



# Latent Diffusion Models

Latents are really just the noised images in pixel space

- No "latent space" so far at least
- But, can add by using an autoencoder



# Text-to-Image Generation + Conditional DMs

Lots of approaches! In particular, for text-to-image generation

- All based on similar principles from multimodal training
- Example: for latent diffusion (Rombach et al '22)
  - “Process  $y$  from various modalities (such as language prompts) we introduce a domain specific encoder ... that projects  $y$  to an intermediate representation ... which is then mapped to the intermediate layers of the UNet via a cross-attention layer “



# Bibliography

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- <https://lilianweng.github.io/posts/2018-10-13-flow-models/>
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**Thank You!**