Deep Learning

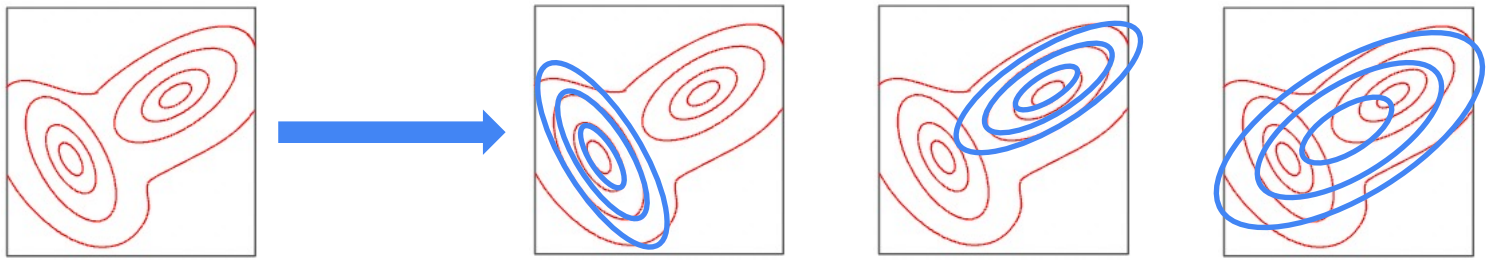
Diffusion Models and Normalizing Flows

11-785 - Fall 2023

Abuzar Khan

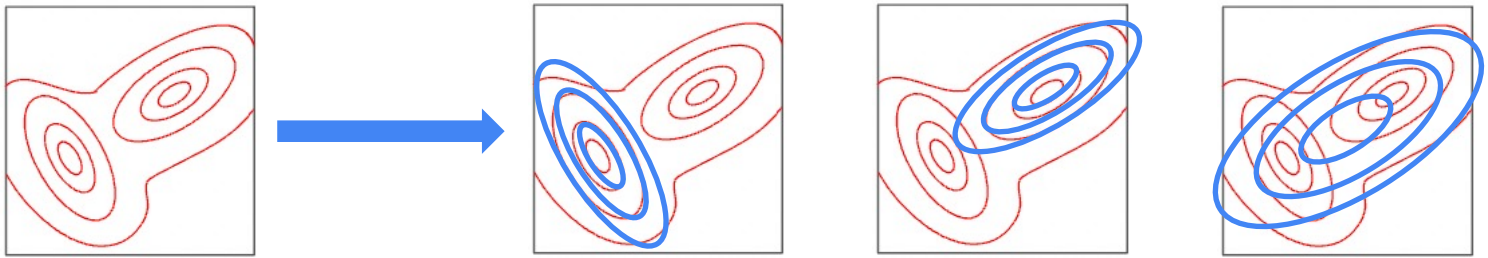
Open Question

1. Say we have a data distribution **p** that is a mixture of two 2D gaussians as shown below in red. We want to approximate this with one gaussian estimate **q** using KL-divergence. Which of the following three will result from optimizing $\mathbf{D}_{\mathrm{KL}}\ (p\ ||\ q)$
2. and which from $\mathbf{D}_{\mathrm{KL}}\ (q\ ||\ p)$?



Open Question

1. Say we have a data distribution **p** that is a mixture of two 2D gaussians as shown below in red. We want to approximate this with one gaussian estimate **q** using KL-divergence. Which of the following three will result from optimizing $\mathbf{D}_{\mathrm{KL}}\ (p\ ||\ q)$
2. and which from $\mathbf{D}_{\mathrm{KL}}\ (q\ ||\ p)$?



$\mathbf{D}_{\mathrm{KL}}\ (q\ ||\ p)$ $\mathbf{D}_{\mathrm{KL}}\ (q\ ||\ p)$ $\mathbf{D}_{\mathrm{KL}}\ (p\ ||\ q)$

Background

1. Generative Models and Discriminative models
2. Autoencoders
3. Variational Autoencoders
    1. Reparameterization trick
    2. ELBO

Sandcastles

How to create a sandcastle:

Step 1: Take a sandcastle

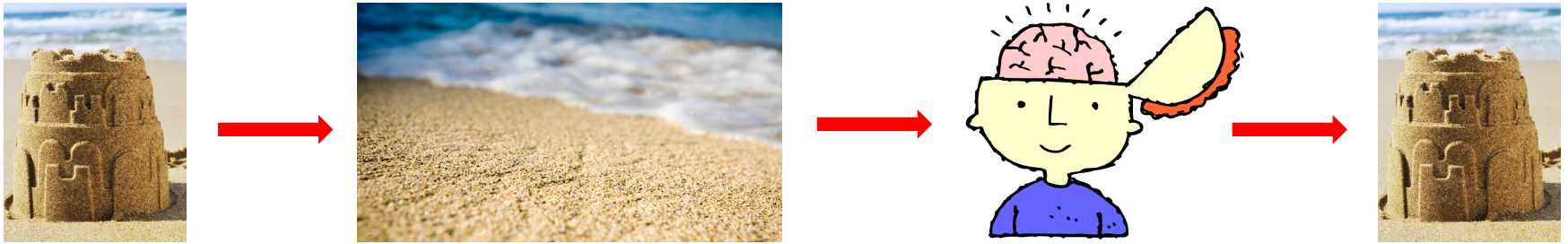
Step 2: Destroy the sandcastle

Step 3: Remember how you destroyed the sandcastle

Step 4: Reverse the process

Key Idea

Once you know how to reconstruct sandcastles, you can start with some different “sand”, apply this process, and end up with a different “sandcastle”



Part 1

Diffusion Models

ELBO Recap

Why use ELBO?

Directly maximizing $p(x)$ is very difficult:

- it involves either marginalizing over the entire latent space $\boldsymbol{Z}$ (intractable for complex models) OR
- It involves having access to the ground truth latent encoder $p(z|x)$

ELBO:

$$\log\big(p(x)\big) \geq \mathbb{E}_{q_\phi(z\,|\,x)}\left[\log \frac{p(x,z)}{q_\phi(z|x)}\right]$$

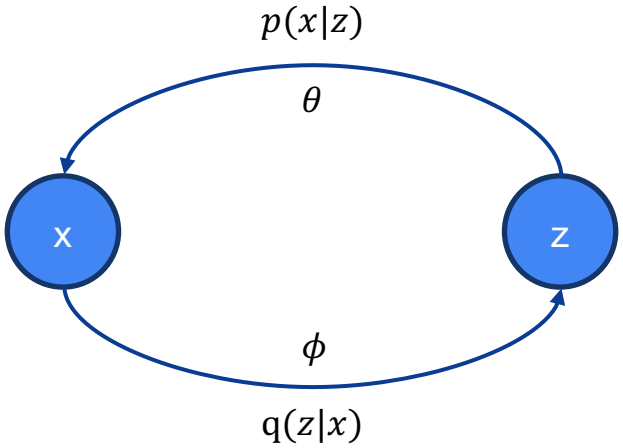
Question: Why does the $\geq$ show up here? → With the derivation in the appendix, we see a $D_{KL}(q_\phi(z|x)\,||\,p(z|x))$ term show up which is always ≥ 0.

Applying chain-rule of probabilities:

$$ELBO = \mathbb{E}_{q_\phi(z\,|\,x)}[\log p_\theta(x|z)] - D_{KL}\big(q_\phi(z\,|\,x)||\,p(z)\big)$$

Reconstruction Prior matching

Variational Autoencoder Recap



Latent variable sampling: $z \sim \mathcal{N}(z;\ \mu_\phi(x), \sigma^2_\phi(x))$

Reparameterization trick: $z = \mu_\phi(x) + \sigma_\phi(x)\ \odot \epsilon,\ \epsilon \sim \mathcal{N}(0, I)$

Training:

- Jointly optimize θ and ϕ
- Maximize **ELBO**

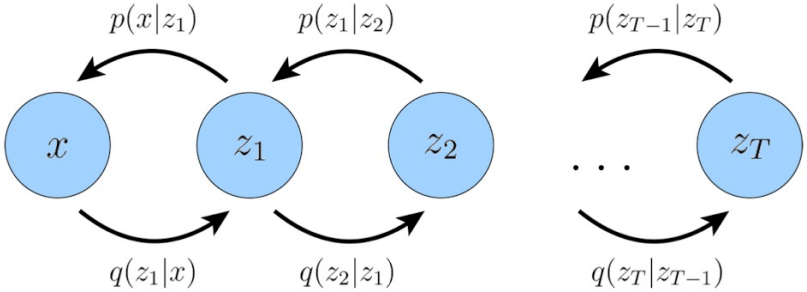
Empirically, we found that two things make VAEs work really well:

1. Increasing the depth of the networks
2. Introducing a hierarchy of latent variables (latent variables of latent variables)

$x \leftarrow z_1 \leftarrow z_2 \leftarrow \ \ldots \leftarrow z_T$, such that each latent is conditioned on all previous latents.

We are particularly interested in such HAVEs that where the process is a **Markovian chain - MHVAE**

Markovian Hierarchical Variational Autoencoder



Joint probability: $p(x, z_{1:T}) = p(z_T)p_\theta(x \mid z_1) \prod_{t=2}^{T} p_\theta(z_{t-1}|z_t)$

Posterior probability: $q_\phi(z_{1:T} \mid x) = q_\phi(z_1| x) \prod_{t=2}^{T} q_\phi(z_t|z_{t-1})$

Updated ELBO:

$$\log\big(p(x)\big) \geq \mathbb{E}_{q_\phi(z_{1:T} \mid x)}\left[\log \frac{p(x, z_{1:T})}{q_\phi(z_{1:T} \mid x)}\right]$$

Diffusion Models

Diffusion models are essentially **MHVAEs** with **3 restrictions:**

1. Latent dimension is the same as the data dimension
2. The encoder has no parameters to be learnt. It is defined to be a linear gaussian such that the t^{th} gaussian is centered around the previous latent z_{t-1}
3. The parameters for the gaussians are scheduled such that the final latent is a standard gaussian.

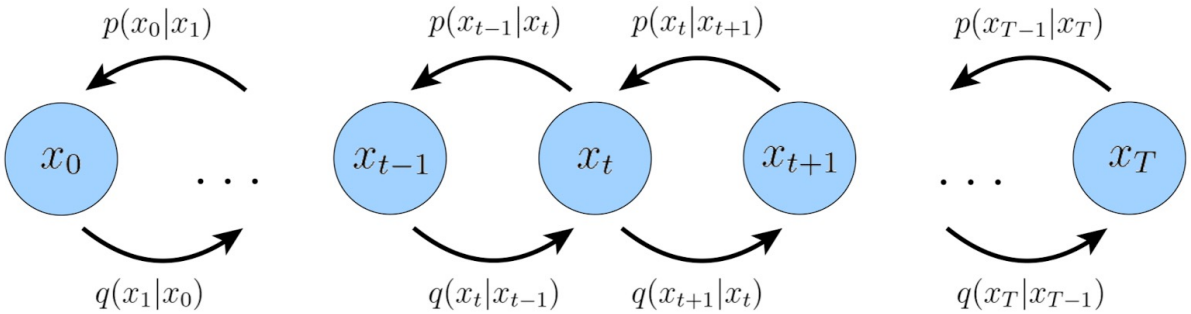
$$z_T \sim \mathcal{N}(z_T;\ 0, \boldsymbol{I})$$

The first restriction allows for some mild abuse of notation:

$$q_\phi(x_{1:T} \mid x_0) = \prod_{t=1}^{T} q_\phi(x_t|x_{t-1})$$

$$p(x_{0:T}) = p(x_T) \prod_{t=1}^{T} p_\theta(x_{t-1}|x_t)$$

Diffusion Models



The first restriction allows for some mild abuse of notation:

$$q_\phi(x_{1:T} \mid x_0) = \prod_{t=1}^{T} q_\phi(x_t|x_{t-1})$$

$$p(x_{0:T}) = p(x_T) \prod_{t=1}^{T} p_\theta(x_{t-1}|x_t)$$

Diffusion Models – Diffusion Process

Following the second restriction, we now define the linear gaussian for the encoding (diffusion) process:

$$q(x_t|x_{t-1}) = \mathcal{N}(x_t; \mu_t(x_{t-1}), \Sigma_t \boldsymbol{I})$$
$$\mu_t(x_{t-1}) = \sqrt{1-\beta_t}x_{t-1}, \quad \Sigma_t = \beta_t$$

We additionally define $\alpha_t = 1 - \beta_t$.
β_t is defined to preserve variance across the diffusion steps.

We can now write

$$q(x_t|x_{t-1}) = \mathcal{N}\left(x_t; \sqrt{\alpha_t}x_{(t-1)}, (1-\alpha_t)\boldsymbol{I}\right)$$

Using the reparameterization trick:

$$x_t = \sqrt{\alpha_t}x_{(t-1)} + \left(\sqrt{1-\alpha_t}\right)\epsilon, \qquad \epsilon \sim \mathcal{N}(0, \boldsymbol{I})$$

$$= \sqrt{\alpha_t\alpha_{t-1}}x_{(t-2)} + \left(\sqrt{1-\alpha_t\alpha_{t-2}}\right)\epsilon$$
$$= \sqrt{\alpha_t\alpha_{t-1}\alpha_{t-2}}x_{(t-3)} + \left(\sqrt{1-\alpha_t\alpha_{t-2}\alpha_{t-3}}\right)\epsilon$$
$$= \sqrt{\alpha_t\alpha_{t-1}\dots\alpha_1}x_{(0)} + \left(\sqrt{1-\alpha_t\alpha_{t-2}\dots\alpha_1}\right)\epsilon$$

$$= \sqrt{\overline{\boldsymbol{\alpha_t}}}\boldsymbol{x_{(0)}} + \left(\sqrt{\boldsymbol{1}-\overline{\boldsymbol{\alpha_t}}}\right)\boldsymbol{\epsilon}$$

This takes us from time step 0 to t in **one step!**

From the third restriction, we get

$$\alpha_T \to 0$$

Sum of two gaussians is another gaussian with mean as the sum of the two means and variance as the sum of the two variances.

$$(1-\alpha_t)\epsilon \to \mathcal{N}(\epsilon; 0, 1-\alpha_t\boldsymbol{I})$$

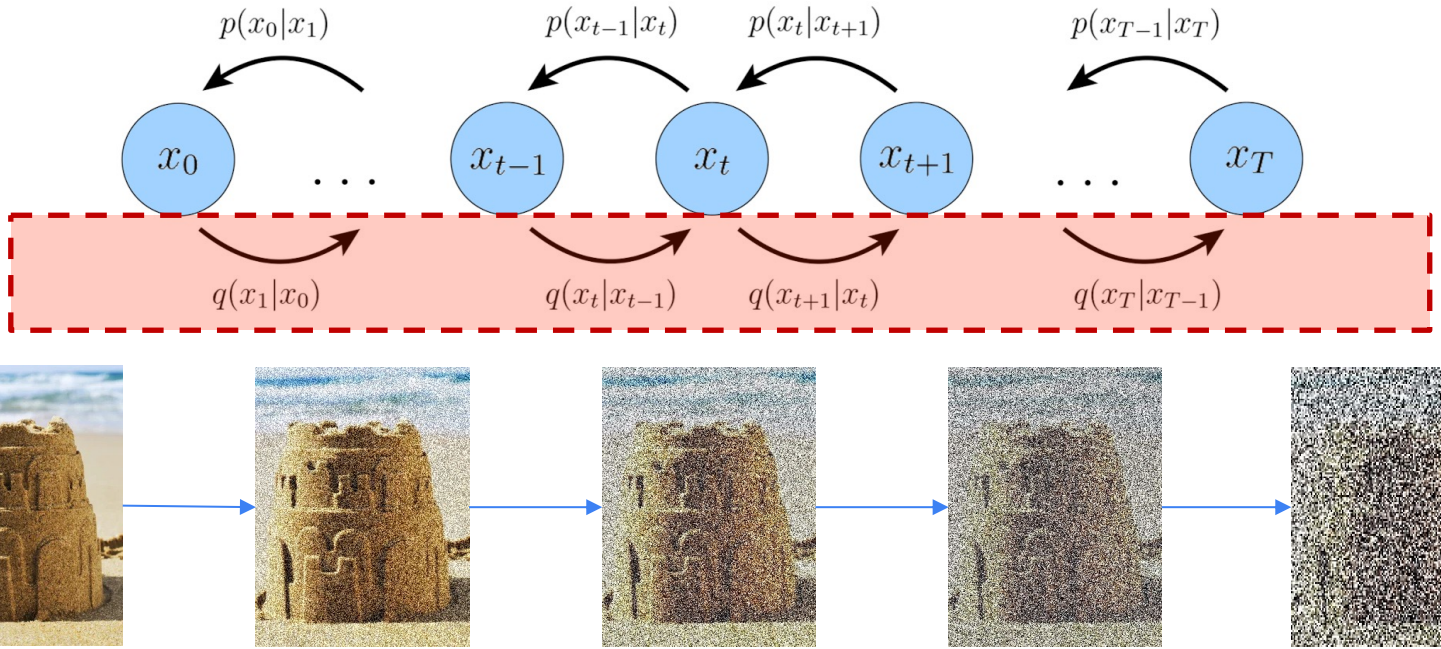
Define

$$\overline{\alpha_t} = \prod_{s=1}^{t}\alpha_s$$

Diffusion Models – Diffusion Process

$$q(x_t|x_{t-1}) = \sqrt{\alpha_t}x_{(t-1)} + \left(\sqrt{1-\alpha_t}\right)\epsilon, \qquad \epsilon \sim \mathcal{N}(0, \boldsymbol{I})$$

This formulation essentially paints a picture of this process to be incrementally adding noise till we reach x_T which is defined to be pure noise.



Diffusion Models – Generative Process

From the third assumption, we can write the exact prior on the final step x_T :

$$p(x_T) = \mathcal{N}(x_T; 0, \boldsymbol{I})$$

For all other steps, we can write a learned distribution:

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

Exactly tractable variance

Neural Network: **U-Net**
Denoising network

$$p_\theta(x_{0:T}) = p(x_T) \prod_{t=1}^{T} p_\theta(x_{t-1}|x_t)$$

Diffusion Models – Updated ELBO

*Derivation in appendix!

$$\log p(x) = \log \int p(x_{0:T}) dx_{0:T}$$

$$\dots *$$

$$= \mathbb{E}_{q(x_1 \mid x_0)}[\log p_\theta(x_0|x_1)] - \mathbb{E}_{q(x_{T-1} \mid x_0)}[D_{KL}(q(x_T \mid x_{T-1})| \,|p(x_T))]$$

$$- \sum_{t=1}^{T-1} \mathbb{E}_{q(x_{t-1}, x_{t+1}| x_0)}\left[D_{KL}\left(q(x_t| x_{t-1}) \,||p_\theta(x_t|x_{t+1})\right)\right]$$

This has **2** random variables for each **t**, this makes the computation slightly hard. We would prefer for there to be need for just **1**!

We can arbitrarily modify the diffusion process distribution to

$$q(x_t| x_{t-1}, x_0) = \frac{q(x_{t-1}| x_t, x_0) q(x_t|x_0)}{q(x_{t-1}| x_0)}$$

Diffusion Models – Updated ELBO

*Derivation in appendix!

$$\log p(x) = \log \int p(x_{0:T}) dx_{0:T}$$

$$\dots *$$

$$= \mathbb{E}_{q(x_1 \mid x_0)}[\log p_\theta(x_0|x_1)] - D_{KL}\left(q(x_T \mid x_0)| \, |p(x_T)\right) - \sum_{t=2}^{T} \mathbb{E}_{q(x_t \mid x_0)}\left[D_{KL}\left(q(x_{t-1}| \, x_t, x_0) \, ||p_\theta(x_{t-1}|x_t)\right)\right]$$

Reconstruction | Prior matching | Denoising

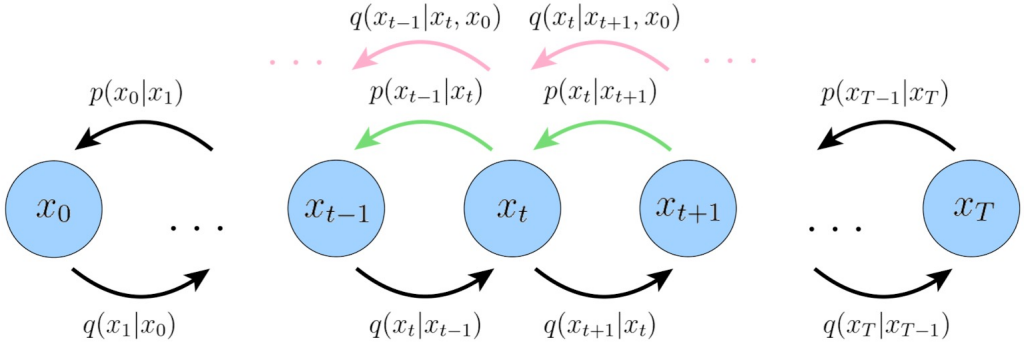
- Reconstruction: Reconstruction from least noisy version **(hyperparameter choice can make this arbitrarily small)**
- Prior matching: Moving the posterior closer to the true prior on the final noisy step **(0 for diffusion models)**
- Denoising: Divergence between approximate denoising (p_θ) and true denoising (q) steps

$q(x_{t-1}| \, x_t, x_0)$ is **tractable** and can be calculated **exactly** without any approximation:

$$q(x_{t-1}| \, x_t, x_0) = \mathcal{N}\left(x_{t-1};\, \bar{\mu}_t, \Sigma_t \boldsymbol{I}\right)$$

$$\bar{\mu}_t = \frac{\sqrt{\alpha_t}(1-\bar{\alpha}_{t-1})x_t + \sqrt{\bar{\alpha}_{t-1}}(1-\alpha_t)x_0}{1-\bar{\alpha}_t}, \qquad \Sigma_t = \frac{(1-\alpha_t)(1-\bar{\alpha}_{t-1})}{1-\bar{\alpha}_t}$$

Diffusion Models – Loss formulation



Loss can focus on the denoising term. Decomposing for each timestep, we can have the t^{th} loss term:

$$L_t = D_{KL}\big(q(x_{t-1}|\, x_t, x_0)\,||p_\theta(x_{t-1}|x_t)\big) + C$$

Since both inputs of the divergence are gaussians, this further simplifies to:

$$\mathrm{L_t} = \mathbb{E}_q\left[\frac{1}{2\Sigma_t}\left|\left|\bar{\mu}_t - \mu_\theta(\mathrm{x}_t, t)\right|\right|^2\right] + \mathrm{C}$$

Diffusion Models – Loss formulation

Further, we have $x_t = \sqrt{\alpha_t} x_{(t-1)} + \left(\sqrt{1-\alpha_t}\right)\epsilon,\ \epsilon \sim \mathcal{N}(0, \boldsymbol{I})$ from definition

This lets us rewrite the true mean of the denoising process as:

$$\bar{\mu}_t = \frac{1}{\sqrt{\alpha_t}}\left(x_t - \frac{(1-\alpha_t)}{\sqrt{1-\bar{\alpha}_t}}\epsilon\right)$$

We can also write the predicted mean as:

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

This lets us reformulate the loss to present a noise prediction problem:

$$L_{t-1} = \mathbb{E}_{x_0,\epsilon}\left[\frac{(1-\alpha_t)^2}{2\Sigma_t \alpha_t (1-\bar{\alpha}_t)}\left|\left|\epsilon - \epsilon_\theta(x_t, t)\right|\right|^2\right] + C$$

Diffusion Models – Training and Inference

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 $\nabla_\theta \left\| \boldsymbol{\epsilon} - \boldsymbol{\epsilon}_\theta(\sqrt{\bar{\alpha}_t}\mathbf{x}_0 + \sqrt{1-\bar{\alpha}_t}\boldsymbol{\epsilon}, t) \right\|^2$ x_t
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** $\mathbf{x}_0$

How do we tell the model what timestep we are on?

Temporal encodings in the form or sinusoids (or anything, really)

Diffusion models - Summary

- Diffusion models are **Markovian Hierarchical VAEs** with extra restrictions
- The loss is the vanilla VAE ELBO loss with an added denoising term
- The encoder has **0 parameters**
- The true denoising posterior can be **exactly calculated**
- The problem can be reformulated as a noise prediction problem
- There’s a ton of math underlying a rather simple intuition

Part 2

Normalizing Flows

Sandcastles

How to create a sandcastle:

Step 1: Take a sandcastle

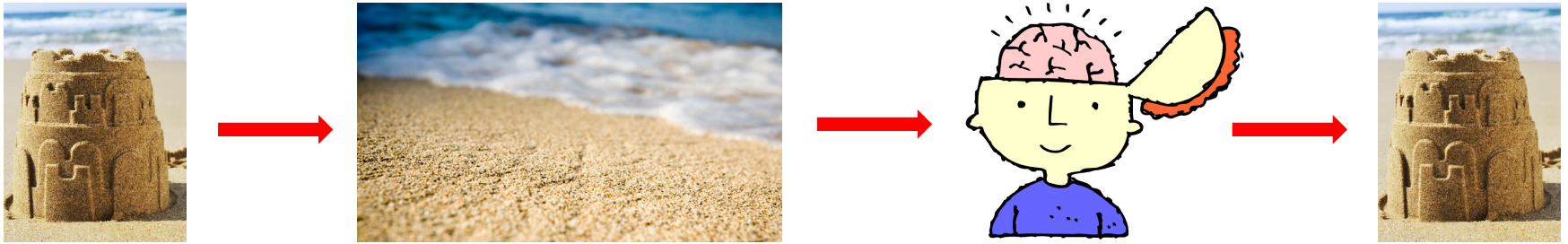
Step 2: Destroy the sandcastle

Step 3: Remember how you destroyed the sandcastle

Step 4: Reverse the process

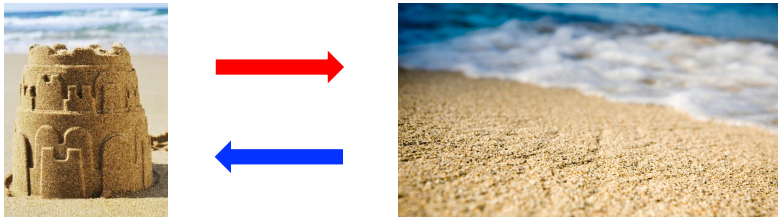
Key Idea

Once you know how to reconstruct sandcastles, you can start with some different “sand”, apply this process, and end up with a different “sandcastle”

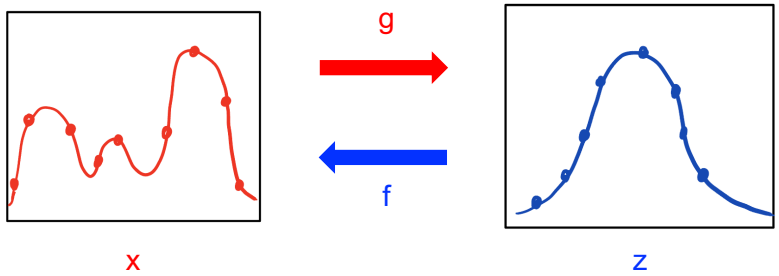


Normalizing Flows – Motivation

- In VAEs we are faced with an intractable likelihood calculation
- We use an ELBO instead as a surrogate objective to MLE
- What if we wanted to do MLE exactly?
- That would require us to go from sandcastle to sand, and back, without any approximation or estimation!

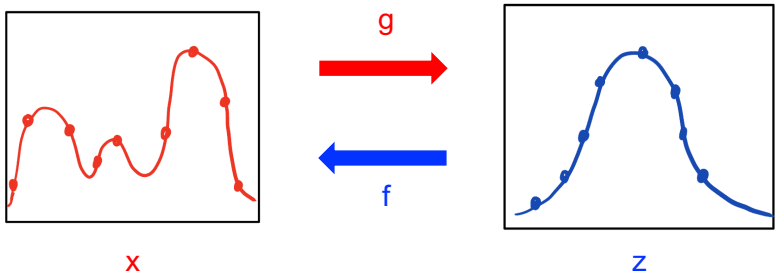


We would need for all the steps we do to be **invertible!**



It follows that
$\mathbf{f^{-1} = g}$

Normalizing Flows – Log likelihood



Bijection (and invertibility) allow us to directly compute the likelihood:

$$\int p_x(x)dx = \int p_z(g(x))dz$$

$$p_x(x) = p_z(g(x))\left|\frac{dg(x)}{dx}\right| \rightarrow p_z(g(x))|det.J(g(x))|$$

$$\log p_x(x) = \log p_z(g(x)) + \log|det.J(g(x))|$$

In multiple dimensions, we generalize to the determinant of the Jacobian

Intuitively

z = g(x) determines where a point in x-space maps to z-space (where to move grains of sand)

|det. J(g(x))| describes how much probability mass (sand) gets moved in a local neighborhood.

Normalizing Flows – Closer look at the Jacobian

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

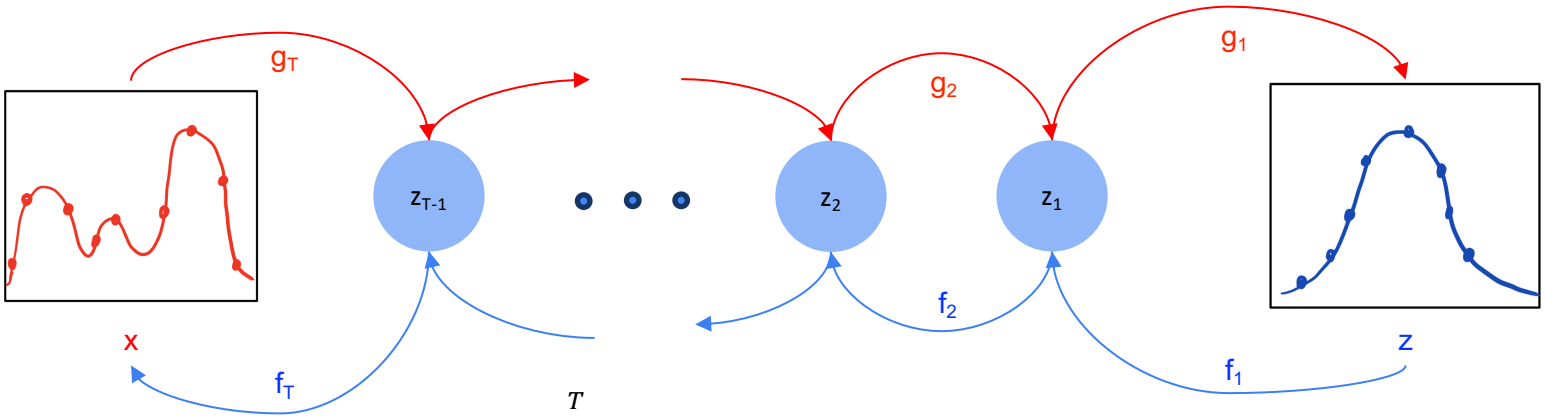
$$J_x g(x) = \begin{bmatrix} \frac{\partial z_1}{\partial x_1} & \cdots & \frac{\partial z_1}{\partial x_K} \\ \vdots & & \vdots \\ \frac{\partial z_K}{\partial x_1} & \cdots & \frac{\partial z_K}{\partial x_K} \end{bmatrix}$$

Normalizing Flows – Compositions

Bijections allow for composing several functions together!

This follows that we can now define:

$$z \sim p(z)$$
$$x = f_T \circ f_{T-1} \circ \cdots \circ f_1\,(z)$$



$$\log p_x(x) = \log p_z(z_0) + \sum_{t=1}^{T} \log\left|det.J_{z_t}\left(g(z_{t-1})\right)\right|, \qquad z_T = x,\ z_0 = z$$

Inverse: $z = g_1 \circ g_2 \circ \cdots \circ g_T(x)$

Normalizing Flows – Characteristics

For a good (efficient) flow, we must have functions (steps) that are:

1. Expressive
2. Invertible
3. Offer cheap to compute Jacobian determinants

Computing a determinant is a **cubic** operation, but some special cases of matrices can make it very cheap.
Especially, **diagonal** matrices:

For a diagonal matrix, the determinant is simply the product of its diagonal elements. Same applies for any triangular matrix!

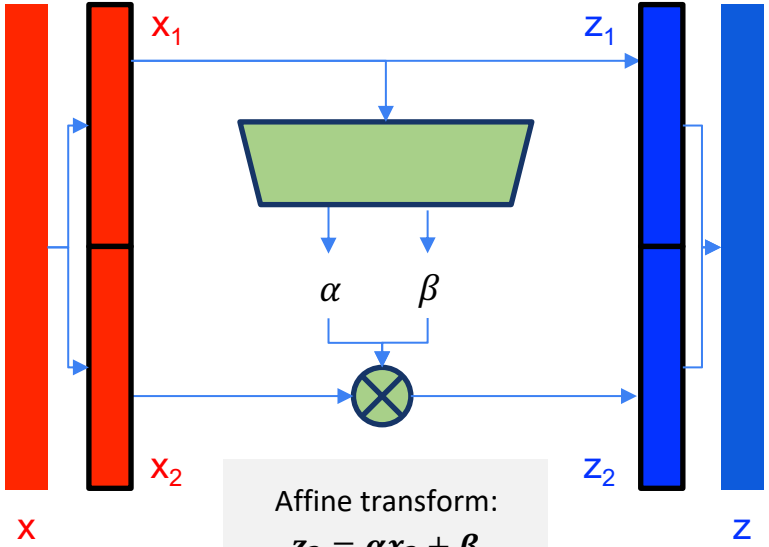
In linear algebra, a **diagonal matrix** is a matrix (usually a square matrix) in which the entries outside the main diagonal (↘) are all zero. The diagonal entries themselves may or may not be zero. Thus, the matrix $D = (d_{ij})$ with n columns and n rows is diagonal if:

$$d_{i,j} = 0 \text{ if } i \neq j \ \forall i, j \in \{1, 2, \ldots, n\}$$

For example, the following matrix is diagonal:

$$\begin{bmatrix} 1 & 0 & 0 \\ 0 & 4 & 0 \\ 0 & 0 & -2 \end{bmatrix}$$

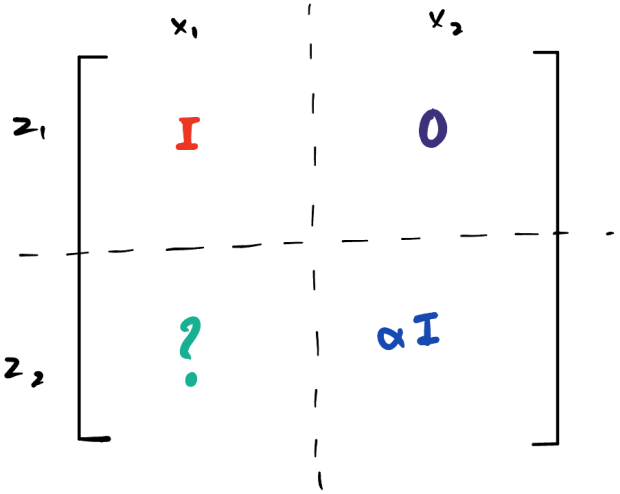
Normalizing Flows – Construction



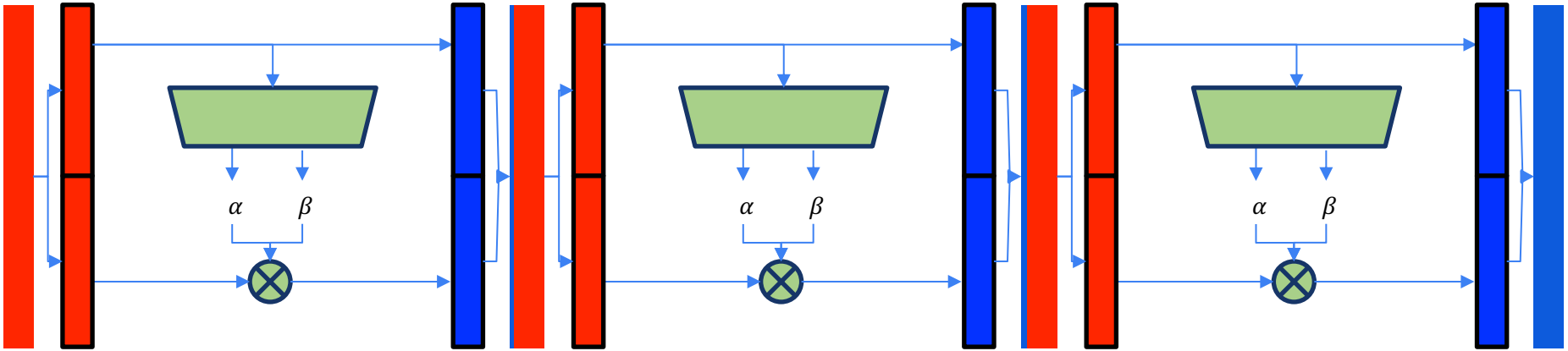
Affine transform:

$$z_2 = \alpha x_2 + \beta$$

$J_a(z)$

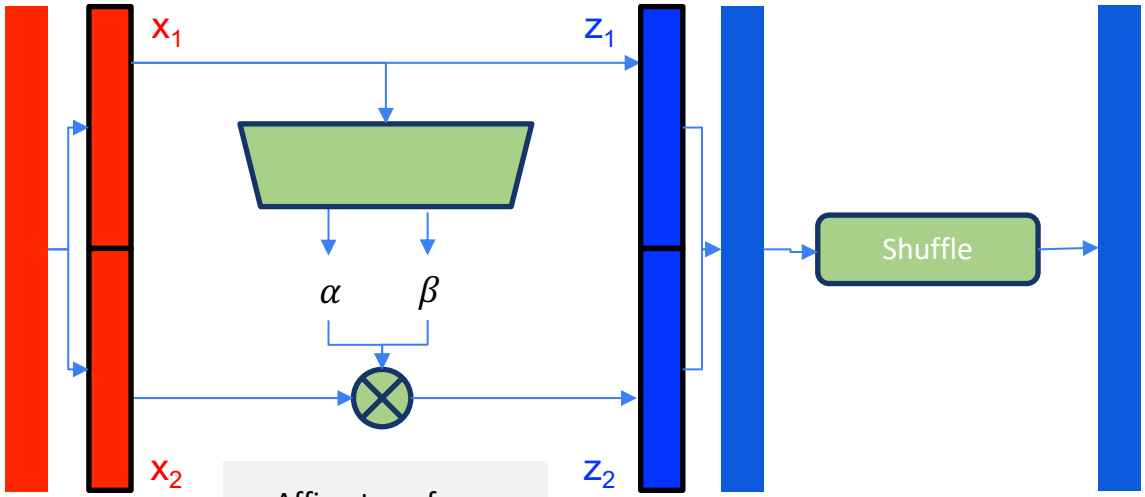


Normalizing Flows – Composition



Is there a problem here?

Normalizing Flows – Construction with a shuffle



Affine transform:

$$z_2 = \alpha x_2 + \beta$$

This the most popular type of flow called as **Coupling Flow** – Used in implementations such as NICE and GLOW

Normalizing Flows – In practice for images

- Multiscale architecture
- Split along channels
- Employ CNNs
- Perform permutations using 1x1 Conv layers

GLOW!

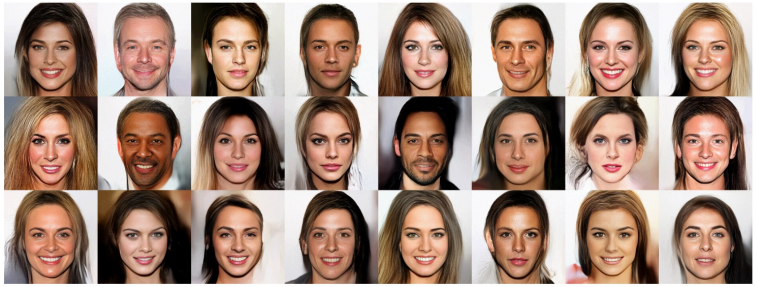


Figure 4: Random samples from the model, with temperature 0.7.

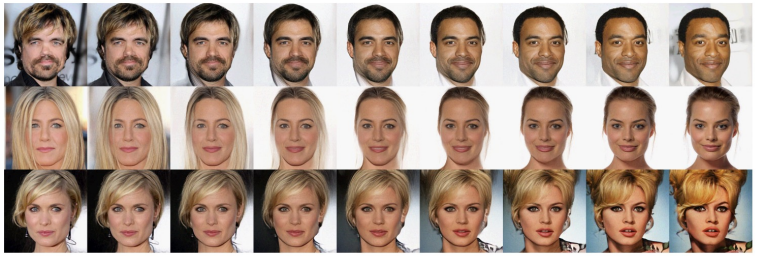


Figure 5: Linear interpolation in latent space between real images.

References

- Denoising Diffusion Probabilistic Models, https://arxiv.org/abs/2006.11239
- Diffusion Models: A Comprehensive Survey of Methods and Applications, https://arxiv.org/abs/2209.00796
- Understanding Diffusion Models: A Unified Perspective, https://arxiv.org/abs/2208.11970
- GLOW, https://proceedings.neurips.cc/paper/2018/file/d139db6a236200b21cc7f752979132d0-Paper.pdf
- Variational inference with normalizing flows, https://arxiv.org/abs/1505.05770
- Normalizing Flows: An Introduction and Review of Current Methods, https://arxiv.org/abs/1908.09257

- Didrik Nielsen's lecture, https://www.youtube.com/watch?v=2tVHbcUP9b8
- Hans van Gorp's lecture, https://www.youtube.com/watch?v=yxVcnuRrKqQ&t=17s
- Tim Salimans' lecture, https://www.youtube.com/watch?v=pea3sH6orMc

Appendix