Diffusion networks

Quick maffs and foundations

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Data generation is a long-standing challenge in machine learning.



Variational Autoencoders (VAEs) are extensively used to project data into a probabilistic latent space from which we can sample to generate new data.



Figure 1: Weng, Lilian. (2018). From Autoencoder to Beta-VAE. https://lilianweng.github.io/posts/2018-08-12-vae/ Generative Adversarial Networks (GANs) have long been used to generate fake data:



Figure 2: From Google advanced machine learning course on GANs. Consulted on Nov. 20 2024.

Outline

2 different perspectives to get somewhat similar conclusions

Score-based generation models

Score matching

Langevin dynamics

Challenges

Multiple noise perturbations

Stochastic differential equations

Denoising diffusion probabilistic models

Forward process

Reverse process

Score-based generation models

Dataset $\{x_1, ..., x_n\}$

Real distribution (unknown): $p_d(x)$

We want to model another distribution $p_m(x; \theta)$ (θ being the model's parameters)

Real distribution (unknown): $p_d(x)$ Modelled distribution: $p_m(x; \theta)$

Find θ such that $p_m(x; \theta)$ is as close as possible to $p_d(x)$

Maximum Likehood:

$$\hat{\theta}_{MLE} = \arg \max_{\theta} \log p_m(x; \theta)$$

Since $p_m(x; \theta)$ is a **normalized** density function, we have:

$$p_m(x;\theta) = rac{ ilde{p}_m(x;\theta)}{Z_{ heta}}$$
 where $Z_{ heta} = \int ilde{p}_m(x;\theta) dx$

 $\tilde{p}_m(x;\theta)$: unnormalized density function Z_{θ} : normalizing constant (intractable because of \int on all data coming from real distribution) Instead of directly maximizing the likelihood, we find a θ such that the gradients of the model's log-likelihood are approx. the same as the gradients of the data distribution log-likelihood.

$$s_{\theta}(x) = \nabla_{x} \log p_{m}(x; \theta)$$
$$= \nabla_{x} \log(\frac{\tilde{p}_{m}(x; \theta)}{Z_{\theta}})$$
$$= \nabla_{x} \log \tilde{p}_{m}(x; \theta) - \nabla_{x} \log Z_{\theta}$$
$$= \nabla_{x} \log \tilde{p}_{m}(x; \theta)$$



To minimize the distance between the score and the score of the real data, we minimize the **Fisher divergence** between the two distributions:

Fisher divergence Doesn't require the two distributions to be normalized.

$$\begin{split} \hat{\theta}_{SM} &= \arg\min_{\theta} D_F(p_d(x) || p_m(x; \theta)) \\ &= \arg\min_{\theta} \frac{1}{2} \mathbb{E}_{p_d(x)} \left[\| \nabla_x \log p_d(x) - \nabla_x \log p_m(x; \theta) \|^2 \right] \\ &= \arg\min_{\theta} \frac{1}{2} \mathbb{E}_{p_d(x)} \left[\| \nabla_x \log p_d(x) - s_\theta(x) \|^2 \right] \end{split}$$

Score matching

However, we have no way of computing the ground truth score $\nabla_x \log p_d(x)$.

$$\hat{\theta}_{SM} = \arg\min_{\theta} \frac{1}{2} \mathbb{E}_{p_d(x)} \left[\left\| \underbrace{\nabla_x \log p_d(x)}_{intractable} - s_{\theta}(x) \right\|^2 \right]$$

There's a derivation from [2] that allows to get to the following objective not involving the ground truth score:

$$\frac{1}{2} \mathbb{E}_{p_d(x)} \left[\left\| \nabla_x \log p_d(x) - s_\theta(x) \right\|^2 \right] \approx \\ \frac{1}{2} \mathbb{E}_{p_d(x)} \left[s_\theta(x)^2 \right] + \mathbb{E}_{p_d(x)} \left[\nabla_x s_\theta(x) \right]$$

We will see a similar derivation later on that is used in practice.

Given that $s_{\theta}(x) = \nabla_x p_m(x_i; \theta)$, we can follow gradient ascent to sample from the model's distribution:

Algorithm 1: Simple sampling

 $\tilde{x}_{0} \sim N(0, I)$ for i = 1..K do $\tilde{x}_{i+1} \leftarrow \tilde{x}_{i} + \alpha \underbrace{\nabla_{x} \log p_{m}(\tilde{x}_{i}; \theta)}_{s_{\theta}(\tilde{x}_{i})}$



Source: Outlier. (Oct 2024). Diffusion Models From Scratch — Score-Based Generative Models Explained. YouTube.

This is problematic as sampling different points will always give the same results (e.g. same image).



Source: Outlier. (Oct 2024). Diffusion Models From Scratch — Score-Based Generative Models Explained. YouTube.

It's almost the same as simple sampling, but with added noise:

 $\frac{\text{Algorithm 2: Langevin dynamics}}{\tilde{x}_{0} \sim \mathcal{N}(0, I)}$ for $i = 1..\mathcal{K}$ do $\epsilon \sim \mathcal{N}(0; 1)$ $\tilde{x}_{i+1} \leftarrow \tilde{x}_{i} + \alpha \underbrace{\nabla_{x} \log p_{m}(\tilde{x}_{i}; \theta)}_{s_{\theta}(\tilde{x}_{i})} + \sqrt{2\alpha\epsilon}$ This enables us to get different data samples. The points will still converge to higher density regions, but with small variations.



Source: Outlier. (Oct 2024). Diffusion Models From Scratch — Score-Based Generative Models Explained. YouTube.

Langevin dynamics

This enables us to get different data samples. The points will still converge to higher density regions, but with small variations.



Source: Song, Yang. (May 2021). Generative Modeling by Estimating Gradients of the Data Distribution. Yang Song's blog.

Naive score-based generative modelling



Figure 3: Score-based generative modeling with Langevin dynamics

Song, Yang. (2021). Generative Modeling by Estimating Gradients of the Data Distribution. https://yang-song.net/blog/2021/score/.

The main pitfall is that the score function is not accurate in low-density regions. When data reside into a high dimensional space, it's very unlikely that we end up in a high-density region.

Figure 4: Estimated scores are only accurate in high density regions.



To bypass having low-density regions, we can add a certain **level** of noise to our data. That create more variability in the data, thus widening the high-density regions to cover more space.



With more space covered by the data, the score function will be accurate in more regions.

However, 2 things to consider:

- 1. Adding too little noise will cause a **inaccurate score function**.
- 2. Adding too much noise will corrupt the data too much.



Multiple noise perturbations

Solution: apply multiple different levels of noise to the data and learn from that **simultaneously**.

We perturb the data distribution $p_d(x)$ with each of the Gaussian noise $\mathcal{N}(0, \sigma_o^2 I)$ with $\sigma_1 < \sigma_2 < \cdots < \sigma_L$. Usually L > 1000.



i-th noise-perturbed distribution:

$$p_{\sigma_i}(x) = \int p(y) \mathcal{N}(x; y, \sigma_i^2 I) dy$$

In reality, we simply sample from $p_{\sigma_i}(x)$ using:

$$x + \sigma_i z$$
 where $z \sim \mathcal{N}(0, I)$

To be able to predict the score on different levels of noise, we condition the score network on the noise level i.

We train a Noise Conditional Score Network (NCSN):

$$s_{\theta}(x, i) \approx \nabla_x \log p_{\sigma_i}(x)$$

With an objective weighting all noise levels:

$$L(\theta) = \sum_{i=1}^{L} \lambda(i) \mathbb{E}_{p_{\sigma_i}(x)} \left[\|\nabla_x \log p_{\sigma_i}(x) - s_{\theta}(x, i)\|^2 \right]$$

- $\lambda(i)$: $\mathbb{R} > 0$. Often $\lambda(i) = \sigma_i^2$.
- Train exactly like before with score matching and this new objective.

Denoising score matching

As before, we still can't compute the ground truth score for any noise level:

$$\sum_{i=1}^{L} \lambda(i) \mathbb{E}_{p_{\sigma_i}(x)} \left[\left\| \underbrace{\nabla_x \log p_{\sigma_i}(x)}_{intractable} - s_{\theta}(x, i) \right\|^2 \right]$$

However, we can derive another form similar to the previous score matching objective:

$$\frac{1}{2} \mathbb{E}_{p_{\sigma_i}(\tilde{x})} \left[\|\nabla_x \log p_{\theta}(\tilde{x}) - s_{\theta}(\tilde{x}, i)\|_2^2 \right] \\= \frac{1}{2} \int p_{\theta}(\tilde{x}) (\nabla_x \log p_{\theta}(\tilde{x}) - s_{\theta}(\tilde{x}))^2 d\tilde{x}$$

$$= \underbrace{\frac{1}{2} \int p_{\theta}(\bar{x}) (\nabla_{\bar{x}} \log p_{\sigma}(\bar{x}))^{2} d\bar{x}}_{\text{don't involve } s_{\theta}(\bar{x})} + \frac{1}{2} \int p_{\theta}(\bar{x}) s_{\theta}(\bar{x})^{2} d\bar{x} - \frac{1}{2} \int p_{\theta}(\bar{x}) 2\nabla_{x} \log p_{\sigma}(\bar{x}) s_{\theta}(\bar{x}) d\bar{x}}$$

$$= \frac{1}{2} \int p_{\theta}(\bar{x}) s_{\theta}(\bar{x})^{2} d\bar{x} - \frac{1}{2} \int p_{\theta}(\bar{x}) 2\nabla_{x} \log p_{\sigma}(\bar{x}) s_{\theta}(\bar{x}) d\bar{x}$$

$$= \frac{1}{2} \int p_{\theta}(\bar{x}) s_{\theta}(\bar{x})^{2} d\bar{x} - \int p_{\theta}(\bar{x}) \frac{\nabla_{x} p_{\sigma}(\bar{x})}{p_{\sigma}(\bar{x})^{2}} s_{\theta}(\bar{x}) d\bar{x}$$

$$= \frac{1}{2} \int p_{\theta}(\bar{x}) s_{\theta}(\bar{x})^{2} d\bar{x} - \iint p(x) \nabla_{x} p_{\sigma}(\bar{x}|x) s_{\theta}(\bar{x}) dx$$

$$= \frac{1}{2} \int p_{\theta}(\bar{x}) s_{\theta}(\bar{x})^{2} d\bar{x} - \iint p(x) \nabla_{x} p_{\sigma}(\bar{x}|x) s_{\theta}(\bar{x}) dx$$

$$= \frac{1}{2} \int p_{\theta}(\bar{x}) s_{\theta}(\bar{x})^{2} d\bar{x} - \iint p(x) \nabla_{x} p_{\sigma}(\bar{x}|x) s_{\theta}(\bar{x}) dxd\bar{x}$$

$$= \frac{1}{2} \sum_{\bar{x}} p_{\theta}(\bar{x}) \left[\| s_{\theta}(\bar{x}) \|_{2}^{2} \right] - \mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\sigma}(\bar{x})} \left[\nabla_{x} \log p_{\sigma}(\bar{x}|x) s_{\theta}(\bar{x}) \right]$$

$$= \frac{1}{2} \mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\theta}(\bar{x})} \left[\| s_{\theta}(\bar{x}) \|_{2}^{2} - 2\nabla_{x} \log p_{\sigma}(\bar{x}|x) s_{\theta}(\bar{x}) \right]$$

$$= \frac{1}{2} \mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\theta}(\bar{x})} \left[\| s_{\theta}(\bar{x}) \|_{2}^{2} - 2\nabla_{x} \log p_{\sigma}(\bar{x}|x) s_{\theta}(\bar{x}) + \| \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} - \| \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} \right]$$

$$= \frac{1}{2} \mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\theta}(\bar{x})} \left[\| s_{\theta}(\bar{x}) - \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} - \| \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} \right]$$

$$= \frac{1}{2} \mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\theta}(\bar{x})} \left[\| s_{\theta}(\bar{x}) - \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} - \| \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} \right]$$

$$= \frac{1}{2} \mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\theta}(\bar{x})} \left[\| s_{\theta}(\bar{x}) - \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} \right] - \underbrace{\mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\theta}(\bar{x})} \left[\| s_{\theta}(\bar{x}) - \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} \right]$$

$$= \frac{1}{2} \mathbb{E}_{x \sim p(\bar{x}), \bar{x} \sim p_{\theta}(\bar{x})} \left[\| s_{\theta}(\bar{x}) - \nabla_{x} \log p_{\sigma}(\bar{x}|x) \|_{2}^{2} \right]$$

$$\frac{1}{2} \mathbb{E}_{x \sim p(x), \tilde{x} \sim p_{\theta}(\tilde{x})} \left[\left\| s_{\theta}(\tilde{x}) - \nabla_{x} \log p_{\sigma}(\tilde{x}|x) \right\|_{2}^{2} \right]$$

 $\tilde{x} = x + \epsilon$ corresponds to $p_{\sigma}(\tilde{x}|x)$. We have:

$$p_{\sigma}(\tilde{x}|x) = \frac{1}{(2\pi)^{d/2}\sigma^2} e^{-\frac{1}{2\sigma^2} \|\bar{x} - x\|^2}$$

$$\nabla_x \log p_{\sigma}(\tilde{x}|x) = \nabla_x \log \left[\frac{1}{(2\pi)^{d/2}\sigma^2} e^{-\frac{1}{2\sigma^2} \|\bar{x} - x\|^2}\right]$$

$$= \nabla_x \log \left[\frac{1}{(2\pi)^{d/2}\sigma^2} + \nabla_x \log \left[e^{-\frac{1}{2\sigma^2} \|\bar{x} - x\|^2}\right]$$

$$= \nabla_x \log \left[e^{-\frac{1}{2\sigma^2} \|\bar{x} - x\|^2}\right]$$

$$= \nabla_x \frac{-1}{2\sigma^2} \|\bar{x} - x\|^2$$

$$= -\frac{2}{2\sigma^2} (\bar{x} - x)$$
Since $\bar{x} = x + \epsilon$

$$= -\frac{1}{\sigma^2} (f + \epsilon - f)$$

$$= -\frac{\epsilon}{\sigma^2}$$

Replacing the objective we found:

$$\frac{1}{2} \mathbb{E}_{x \sim p(x), \tilde{x} \sim p_{\theta}(\tilde{x})} \left[\left\| s_{\theta}(\tilde{x}) - \underbrace{\nabla_{x} \log p_{\sigma}(\tilde{x}|x)}_{2} \right\|_{2}^{2} \right]$$
$$= \frac{1}{2} \mathbb{E}_{x \sim p(x), \tilde{x} \sim p_{\theta}(\tilde{x})} \left[\left\| s_{\theta}(\tilde{x}) + \frac{\epsilon}{\sigma^{2}} \right\|_{2}^{2} \right]$$

Making it **tractable** where $s_{\theta}(\tilde{x})$ is trying to predict $-\frac{\epsilon}{\sigma^2}$ (to remove the noise)!

Annealed Langevin Dynamics

Sampling with NCSN: Same method as Langevin dynamics, but with a twist: **each step**, the scores are predicted from a different noise level.

 Algorithm 3: Annealed Langevin dynamics

 $x_0 \sim \mathcal{N}(0, I)$

 for i = L..1 do

 $\ell \in \mathcal{N}(0; 1)$
 $z_t \sim \mathcal{N}(0; 1)$
 $\alpha_i \leftarrow \epsilon \frac{\sigma_i^2}{\sigma_L^2}$
 $\tilde{x}_{i+1} \leftarrow \tilde{x}_i + \alpha_i \underbrace{\nabla_x \log p_m(x_i; \sigma_i, \theta)}_{s_{\theta}(x_i, i)} + \sqrt{2\alpha_i} z_t$

Important: noise is decreasing at each iteration (since $\sigma_L > \sigma_{L-1} > \cdots > \sigma_1$)

Annealed Langevin Dynamics



Annealed Langevin Dynamics



With NCSN, we have a **discrete and finite sequence of noise levels** that the model denoise data with. Most of the time, the sequence length L is fixed to at least 1000 steps.



In general, an Ordinary Differential Equation (ODE) has the following formulation:

$$dx = f(x, t)dt$$

which describes the evolution of a deterministic system over time.

Now, if the process is random, we can analyse it's evolution with Stochastic differential equations (SDEs):

 $dx = f(x_t, t)dt + g(t)dw_t$

where

- *dw_t* is infinitesimal noise.
- $f(x_t, t)$ is the drift coefficient (deterministic part).
The forward SDE perturbs the data on a continuous time-scale:

$$dx = f(x_t, t)dt + g(t)dw_t$$



It turns out that if you do the derivation, the reverse process of an SDE in general is the following:

$$dx = \left[f(x_t, t) - g^2(t)\nabla_x \log p_t(x)\right] dt + g(t)dw$$
$$= \left[f(x_t, t) - g^2(t)s_\theta(x, t)\right] dt + g(t)dw$$

We just have to train a **Time-dependent score-based model** $s_{\theta}(x, t) \approx \nabla_x \log p_t(x)$ which is basically the same as the NCSN $s_{\theta}(x, i) \approx \nabla_x \log p_{\sigma_i}(x)$. Instead of conditioning on a specific schedule time-step, we **condition on the continuous time** t.

The reverse SDE starts from noise from a prior distribution (gaussian) to the data:

$$dx = \left[f(x_t, t) - g^2(t)s_\theta(x, t)\right]dt + g(t)dw$$



Stochastic differential equations



To train that **Time-dependent score-based model**, we use the following objective:

$$\mathbb{E}_{t \in \mathcal{U}(0,T)} \mathbb{E}_{\rho_t(x)} \left[\lambda(t) \left\| \nabla_x \log p_t(x) - s_{\theta}(x,t) \right\|_2^2 \right]$$

- $\mathcal{U}(0, T)$ a uniform distribution over timesteps
- Weighting function $\lambda(t) \propto rac{1}{\mathbb{E}\left[\left\|
 abla_{x(t)} \log p(x(t) | x(0)) \right\|_2^2\right]}$

We now have the reverse SDE:

$$dx = \left[f(x_t, t) - g^2(t)s_\theta(x, t)
ight]dt + g(t)dw$$

How do we solve it? (Solving means to get from t = T (noise) to t = 0 (data))

- 1. Train the time-dependent score-based model $s_{\theta}(x, t)$
- Use any numerical SDE solver to solve the SDE (e.g. Euler-Maruyama).

The Euler-Maruyama discretizes the SDE into discrete time-steps (similar to Langevin dynamics).

Algorithm 4: Example SDE solver: Euler-Maruyama

$$\begin{split} \Delta t &\approx 0 \text{ (very small)}; \\ t &\leftarrow T; \\ z_t &\sim \mathcal{N}(0, I); \\ \textbf{repeat} \\ & \left| \begin{array}{c} \Delta x \leftarrow \left[f(x, t) - g^2(t) s_\theta(x, t) \right] \Delta t + g(t) \sqrt{|\Delta t|} z_t; \\ x \leftarrow x + \Delta x; \\ t \leftarrow t + \Delta t; \\ \textbf{until } t &\approx 0; \\ \end{split} \right. \end{split}$$

Benefits:

- We can solve using an **arbitrary number of denoising steps** with the same model.
- The SDE solver is independent of the model.



This formulation of score-based models bring:

- Exact likelihood computation (with probability flow ODE).
- Better sampling methods (but still much slower than GANs).
- Link the score-based models to DDPM.



Denoising diffusion probabilistic models

Diffusion models at high-level





The real $q(x_{t-1}|x_t)$ reverse function is unknown. We approximate it with a learned function $p_{\theta}(x_{t-1}|x_t)$. The forward process gradually adds gaussian noise to the data **during training**.

- 1. Sample a data point x_0 from the real data $x_0 \sim q_d(x)$
- 2. Add gaussian noise with variance β_t to x_{t-1} to produce a **new** latent variable $q(x_t|x_{t-1})$ where

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

Problem: Calculating x_t from x_0 iteratively is expensive (especially if it's an image with t > 1000). **Solution**: Reparametrization trick!

From $q(x_t|x_{t-1}) = \mathcal{N}(x_t; \sqrt{1 - \beta_t}x_{t-1}, \beta_t I)$ to $z = \mu + \sigma \odot \epsilon$

Forward process

Solution: Reparametrization trick! Let $\alpha_t = 1 - \beta_t$ and $\overline{\alpha}_t = \prod_{t=1}^T \alpha_t$ $x_t = \mathcal{N}(x_t; \sqrt{1 - \beta_t x_{t-1}, \beta_t I})$ $=\sqrt{1-\beta_t}x_{t-1}+\sqrt{\beta_t}\epsilon_{t-1}$ $= \sqrt{\alpha_t} X_{t-1} + \sqrt{1-\alpha_t} \epsilon_{t-1}$ $= \sqrt{\alpha_t} \left(\sqrt{\alpha_{t-1}} x_{t-2} + \sqrt{1 - \alpha_{t-1}} \epsilon_{t-2} \right) + \sqrt{1 - \alpha_t} \epsilon_{t-1}$ $= \sqrt{\alpha_t \alpha_{t-1}} x_{t-2} + \sqrt{\alpha_t (1 - \alpha_{t-1})} \epsilon_{t-2} + \sqrt{1 - \alpha_t} \epsilon_{t-1}$ $=\sqrt{\alpha_t \alpha_{t-1}} x_{t-2} + \sqrt{1-\alpha_t \alpha_{t-1}} \overline{\epsilon}_{t-2}$ = ... replace x_{t-2} and so on ... $=\sqrt{\overline{\alpha_t}}x_0+\sqrt{1-\overline{\alpha_t}}\epsilon_t$ $q(x_t|x_0) = \mathcal{N}(x_t; \sqrt{\overline{\alpha}_t}x_0, (1-\overline{\alpha}_t)/)$

- Recall: $\alpha_t = 1 \beta_t$ and β_t is an hyperparameter.
- We can precompute $\overline{\alpha}_t$ (not expensive).
- With q(xt|x0) = N(xt; √at x0, (1 − at)I) we can sample xt from x0 in one step.

How to parametrize β_t ?

The choice of the noise schedule can be arbitrary as long as there's a near-linear change in the middle and very subtle changes around t=0 and t=T.

- Linear schedule (perturbs image too quickly initially). E.g. from $\beta_1 = 10^{-4}$ to $\beta_T = 0.02$.
- Cosine schedule (showed better results)

$$\beta_t = \operatorname{clip}(1 - \frac{\overline{\alpha}_t}{\overline{\alpha}_{t-1}}, 0.999) \text{ with } \overline{\alpha}_t = \frac{f(t)}{f(0)}$$

and $f(t) = \cos\left(\frac{t/T + s}{1 + s}\frac{\pi}{2}\right)^2$



Reverse process

If we could sample from $q(x_{t-1}|x_t)$ (intractable), we could recreate an image from a randomly sampled x_T from a gaussian.

If β_t is small enough, $q(x_{t-1}|x_t)$ is gaussian. So we can approximate it with:

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

- $\mu_{\theta}(x_t, t)$ is learned.
- In practice (for DDPM), Σ_θ(x_t, t) = σ_t² l which comes from a fixed variance schedule (not learned).

Objective: minimize $-\log p_{\theta}(x_0)$ **Issue**: $p_{\theta}(x_0)$ is intractable as it depends on all previous timesteps $x_0, x_1, ..., x_{T-1}, x_T$.

Instead of minimizing $-\log p_{\theta}(x_0)$, we minimize the variational lower bound:

 $-\log(p_{\theta}(x_{0})) \leq -\log p_{\theta}(x_{0}) + D_{KL}(q(x_{1:T}|x_{0})||p_{\theta}(x_{1:T}|x_{0}))$

- We add the KL (always ≥ 0) since we are minimizing.
- Still depends on $-\log p_{\theta}(x_0)$! We have to reformulate.

Variational Lower Bound

$$= -\log(p_{\theta}(x_{0}))$$

$$\leq -\log p_{\theta}(x_{0}) + D_{KL}(q(x_{1:T}|x_{0})||p_{\theta}(x_{1:T}|x_{0}))$$

$$= -\log p_{\theta}(x_{0}) + \mathbb{E}_{x_{1:T}\sim q(x_{1:T}|x_{0})} \left[\log \frac{q(x_{1:T}|x_{0})}{p_{\theta}(x_{1:T}|x_{0})}\right] \quad \text{(Def. KL)}$$

$$= -\log p_{\theta}(x_{0}) + \mathbb{E}_{q} \left[\log \frac{q(x_{1:T}|x_{0})}{\frac{p_{\theta}(x_{0:T})}{p_{\theta}(x_{0})}}\right] \quad \text{(Bayes)}$$

$$= -\log p_{\theta}(x_{0}) + \mathbb{E}_{q} \left[\log \frac{q(x_{1:T}|x_{0})}{p_{\theta}(x_{0:T})} + \log p_{\theta}(x_{0})\right]$$

$$L_{VLB} = \mathbb{E}_{q} \left[\log \frac{q(x_{1:T}|x_{0})}{p_{\theta}(x_{0:T})}\right]$$

 $q(x_{1:T}|x_0)$ is the forward process, we can compute it. $p_{\theta}(x_{0:T})$ is not analytically computable.

Variational Lower Bound

We have to reformulate L_{VLB} even more to make it computable:

$$\begin{split} \mathcal{L}_{VLB} &= \mathbb{E}_{q} \left[\log \frac{q(x_{1:T} \mid x_{0})}{p_{\theta}(x_{0:T})} \right] \\ &= \mathbb{E}_{q} \left[\log \frac{\prod_{t=1}^{T} q(x_{t} \mid x_{t-1})}{p_{\theta}(x_{T}) \prod_{t=1}^{T} p_{\theta}(x_{t-1} \mid x_{t})} \right] & \text{(Def. } q \text{ and} \\ &= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=1}^{T} \log \frac{q(x_{t} \mid x_{t-1})}{p_{\theta}(x_{t-1} \mid x_{t})} \right] & \text{(Log prop.)} \\ &= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t} \mid x_{t-1})}{p_{\theta}(x_{t-1} \mid x_{t})} + \log \frac{q(x_{1} \mid x_{0})}{p_{\theta}(x_{0} \mid x_{1})} \right] & \text{(First term)} \end{split}$$

Details

56

$$\begin{split} &= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t}|x_{t-1})}{p_{\theta}(x_{t-1}|x_{t})} + \log \frac{q(x_{1}|x_{0})}{p_{\theta}(x_{0}|x_{1})} \right] \\ &= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t},x_{0})}{p_{\theta}(x_{t-1}|x_{t})} \frac{q(x_{t}|x_{0})}{q(x_{t-1}|x_{0})} + \log \frac{q(x_{1}|x_{0})}{p_{\theta}(x_{0}|x_{1})} \right] \\ &= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t},x_{0})}{p_{\theta}(x_{t-1}|x_{t})} + \sum_{t=2}^{T} \log \frac{q(x_{t}|x_{0})}{q(x_{t-1}|x_{0})} + \log \frac{q(x_{1}|x_{0})}{p_{\theta}(x_{0}|x_{1})} \right] \end{split}$$

Additional simplification

$$\sum_{t=2}^{T} \log \frac{q(x_t|x_0)}{q(x_{t-1}|x_0)} = \log \prod_{t=2}^{T} \frac{q(x_t|x_0)}{q(x_{t-1}|x_0)} = \log \frac{q(x_2|x_0) \ q(x_3|x_0) \ q(x_4|x_0)...}{q(x_1|x_0) \ q(x_2|x_0) \ q(x_3|x_0)...}$$
$$= \log \frac{q(x_T|x_0)}{q(x_1|x_0)}$$

$$= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} + \sum_{t=2}^{T} \log \frac{q(x_{t}|x_{0})}{q(x_{t-1}|x_{0})} + \log \frac{q(x_{1}|x_{0})}{p_{\theta}(x_{0}|x_{1})} \right]$$

$$= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} + \log \frac{q(x_{T}|x_{0})}{q(x_{1}|x_{0})} + \log \frac{q(x_{1}|x_{0})}{p_{\theta}(x_{0}|x_{1})} \right]$$

$$= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} + \log \frac{q(x_{T}|x_{0}) \cdot q(x_{T}|x_{0})}{p_{\theta}(x_{0}|x_{1}) \cdot q(x_{T}|x_{0})} \right]$$

$$= \mathbb{E}_{q} \left[-\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} + \log \frac{q(x_{T}|x_{0})}{p_{\theta}(x_{0}|x_{1})} \right]$$

$$= \mathbb{E}_{q} \left[\log p_{\theta}(x_{T}) + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} + \log q(x_{T}|x_{0}) - \log p_{\theta}(x_{0}|x_{1}) \right]$$

$$= \mathbb{E}_{q} \left[\log \frac{q(x_{T}|x_{0})}{p_{\theta}(x_{T})} + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} - \log p_{\theta}(x_{0}|x_{1}) \right]$$

$$= \mathbb{E}_{q} \left[\log \frac{q(x_{T}|x_{0})}{p_{\theta}(x_{T})} + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} - \log p_{\theta}(x_{0}|x_{1}) \right]$$

$$= \mathbb{E}_{q} \left[\log \frac{q(x_{T}|x_{0})}{p_{\theta}(x_{T})} + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} - \log p_{\theta}(x_{0}|x_{1}) \right]$$

$$= \mathbb{E}_{q} \left[\log \frac{q(x_{T}|x_{0})}{p_{\theta}(x_{T})} + \sum_{t=2}^{T} \log \frac{q(x_{t-1}|x_{t}, x_{0})}{p_{\theta}(x_{t-1}|x_{t})} - \log p_{\theta}(x_{0}|x_{1}) \right]$$

(1) $D_{KL}(q(x_T|x_0)||p_{\theta}(x_T))$

$D_{KL}(q(x_T|x_0)||p_{\theta}(x_T))$

- q(x_T | x₀): no learnable parameter AND converges to isotropic gaussian.
- $p_{\theta}(x_{T})$: random noise sampled from an isotropic gaussian.

We can ignore this term (1)!

We are left with:

$$L_{VLB} = \mathbb{E}_{q}[\underbrace{\sum_{t=2}^{T} D_{KL}(q(x_{t-1}|x_{t}, x_{0})||p_{\theta}(x_{t-1}|x_{t}))}_{(2)} - \underbrace{\log p_{\theta}(x_{0}|x_{1})]}_{(3)}]$$

(2)
$$\sum_{t=2}^{T} D_{KL}(q(x_{t-1}|x_t, x_0)||p_{\theta}(x_{t-1}|x_t))$$

Recall that estimating $q(x_{t-1}|x_t)$ is intractable. However, by conditioning on x_0 , the term in the above equation can be computed:

Let
$$q(x_{t-1}|x_t, x_0) = \mathcal{N}(x_{t-1}; \tilde{\mu}(x_t, x_0), \tilde{\beta}I)$$

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

Replace each term by their gaussian form w/o the coeffcient.

$$\propto \exp\left[\frac{-1}{2}\left(\frac{(x_t - \sqrt{\alpha_t}x_{t-1})^2}{\sqrt{\beta_t}^2} + \frac{(x_{t-1} - \sqrt{\overline{\alpha}_{t-1}x_0})^2}{\sqrt{1 - \overline{\alpha}_{t-1}}^2} - \frac{(x_t - \sqrt{\overline{\alpha}_t}x_0)^2}{\sqrt{1 - \overline{\alpha}_t}^2}\right)\right]$$

$$= \exp\left[\frac{-1}{2}\left(\frac{(x_t - \sqrt{\alpha_t}x_{t-1})^2}{\beta_t} + \frac{(x_{t-1} - \sqrt{\overline{\alpha}_{t-1}x_0})^2}{1 - \overline{\alpha}_{t-1}} - \frac{(x_t - \sqrt{\overline{\alpha}_t}x_0)^2}{1 - \overline{\alpha}}\right)\right]$$

$$\begin{split} &= \exp\left[\frac{-1}{2}\left(\frac{(x_t - \sqrt{\alpha_t}x_{t-1})^2}{\beta_t} + \frac{(x_{t-1} - \sqrt{\overline{\alpha_{t-1}}x_0})^2}{1 - \overline{\alpha_{t-1}}} - \frac{(x_t - \sqrt{\overline{\alpha_t}}x_0)^2}{1 - \overline{\alpha}}\right)\right] \\ &= \exp\left[\frac{-1}{2}\left(\frac{x_t^2 - 2\sqrt{\alpha_t}x_tx_{t-1} + \alpha_tx_{t-1}^2}{\beta_t} + \frac{x_{t-1}^2 - 2\sqrt{\overline{\alpha_{t-1}}x_{t-1}x_0} + \overline{\alpha_{t-1}}x_0^2}{1 - \overline{\alpha_{t-1}}} - \frac{(x_t - \sqrt{\overline{\alpha_t}}x_0)^2}{1 - \overline{\alpha}}\right)\right] \\ &= \exp\left[\frac{-1}{2}\left(\frac{x_t^2}{\beta_t} - \frac{2\sqrt{\alpha_t}x_tx_{t-1}}{\beta_t} + \frac{\alpha_tx_{t-1}^2}{\beta_t} + \frac{x_{t-1}^2}{1 - \overline{\alpha_{t-1}}} - \frac{2\sqrt{\overline{\alpha_{t-1}}x_{t-1}x_0}}{1 - \overline{\alpha_{t-1}}} + \frac{\overline{\alpha_t}x_0^2}{1 - \overline{\alpha}}\right)\right] \end{split}$$

Discard/regroup terms in black that don't depend on x_{t-1}

$$= \exp\left[\frac{-1}{2}\left(-\frac{2\sqrt{\alpha_t}x_tx_{t-1}}{\beta_t} + \frac{\alpha_tx_{t-1}^2}{\beta_t} + \frac{x_{t-1}^2}{1 - \overline{\alpha}_{t-1}} - \frac{2\sqrt{\overline{\alpha_{t-1}}x_{t-1}}x_0}{1 - \overline{\alpha}_{t-1}} + C(x_t, x_0)\right)\right]$$

We factorize x_{t-1}^2 and x_{t-1}
$$= \exp\left[\frac{-1}{2}\left(\frac{\alpha_t}{\beta_t} + \frac{1}{1 - \overline{\alpha}_{t-1}}\right)x_{t-1}^2 - 2\left(\frac{\sqrt{\alpha_t}x_t}{\beta_t} + \frac{\sqrt{\overline{\alpha}_{t-1}}x_0}{1 - \overline{\alpha}_{t-1}}\right)x_{t-1} + C(x_t, x_0)\right]$$

$$\exp\left[\frac{-1}{2}\left(\frac{\alpha_t}{\beta_t}+\frac{1}{1-\overline{\alpha}_{t-1}}\right)x_{t-1}^2-2\left(\frac{\sqrt{\alpha_t}x_t}{\beta_t}+\frac{\sqrt{\overline{\alpha}_{t-1}}x_0}{1-\overline{\alpha}_{t-1}}\right)x_{t-1}+C(x_t,x_0)\right]$$

Following the gaussian distribution form (omitting coefs) of $\mathcal{N}(x_{t-1}, \tilde{\mu}(x_t, x_0), \tilde{\beta}I)$

$$\begin{split} \mathcal{N}(x_{t-1}, \tilde{\mu}(x_t, x_0), \tilde{\beta}I) &= \exp\left[\frac{-1}{2} \frac{(x_{t-1} - \tilde{\mu}(x_t, x_0))^2}{\tilde{\beta}}\right] \\ &= \exp\left[\frac{-1}{2} \frac{x_{t-1}^2 - 2x_{t-1}\tilde{\mu}(x_t, x_0) + \tilde{\mu}(x_t, x_0)^2}{\tilde{\beta}}\right] \\ &= \exp\left[\frac{-1}{2} \left(\frac{x_{t-1}^2}{\tilde{\beta}} - \frac{2x_{t-1}\tilde{\mu}(x_t, x_0)}{\tilde{\beta}} + C(x_t, x_0)\right)\right] \\ &= \exp\left[\frac{-1}{2} \left(\frac{1}{\tilde{\beta}} x_{t-1}^2 - \frac{2\tilde{\mu}(x_t, x_0)}{\tilde{\beta}} x_{t-1} + C(x_t, x_0)\right)\right] \end{split}$$

By plugging the formulation at the top into the gaussian form, we can find $\tilde{\mu}(x_t, x_0)$ and $\tilde{\beta}$.

Recall that $\beta_t = 1 - \alpha_t$. We begin by finding $\tilde{\beta}$:

$$\begin{split} \vec{\beta} &= \left(\frac{\alpha_t}{\beta_t} + \frac{1}{1 - \overline{\alpha}_{t-1}}\right)\\ \vec{\beta} &= \frac{1}{\frac{\alpha_t}{\beta_t} + \frac{1}{1 - \overline{\alpha}_{t-1}}}\\ &= \frac{1}{\frac{\alpha_t - \alpha_t \overline{\alpha}_{t-1} + \beta_t}{\beta_t (1 - \overline{\alpha}_{t-1})}}\\ &= \frac{\beta_t (1 - \overline{\alpha}_{t-1})}{\alpha_t - \alpha_t \overline{\alpha}_{t-1} + \beta_t}\\ &= \frac{1 - \overline{\alpha}_{t-1}}{\alpha_t - \overline{\alpha}_t + \beta_t} \beta_t\\ &= \frac{1 - \overline{\alpha}_{t-1}}{\varphi_t - \overline{\alpha}_t + 1 - \varphi_t} \beta_t\\ &= \frac{1 - \overline{\alpha}_{t-1}}{1 - \overline{\alpha}_t} \beta_t \end{split}$$

Now, for $\tilde{\mu}(x_t, x_0)$, we have:

$$\frac{2\tilde{\mu}(\mathbf{x}_{t}, \mathbf{x}_{0})}{\tilde{\beta}_{t}} = 2\left(\frac{\sqrt{\alpha_{t}}}{\beta_{t}}\mathbf{x}_{t} + \frac{\sqrt{\overline{\alpha}_{t-1}}}{1 - \overline{\alpha}_{t-1}}\mathbf{x}_{0}\right)$$
$$\tilde{\mu}(\mathbf{x}_{t}, \mathbf{x}_{0}) = \left(\frac{\sqrt{\alpha_{t}}}{\beta_{t}}\mathbf{x}_{t} + \frac{\sqrt{\overline{\alpha}_{t-1}}}{1 - \overline{\alpha}_{t-1}}\mathbf{x}_{0}\right)\tilde{\beta}_{t}$$
$$= \left(\frac{\sqrt{\alpha_{t}}}{\beta_{t}}\mathbf{x}_{t} + \frac{\sqrt{\overline{\alpha}_{t-1}}}{1 - \overline{\alpha}_{t-1}}\mathbf{x}_{0}\right)\frac{1 - \overline{\alpha}_{t-1}}{1 - \overline{\alpha}_{t}}\beta_{t}$$
$$= \frac{\sqrt{\alpha_{t}}(1 - \overline{\alpha}_{t-1})}{(1 - \overline{\alpha}_{t})}\mathbf{x}_{t} + \frac{\sqrt{\overline{\alpha}_{t-1}}\beta_{t}}{(1 - \overline{\alpha}_{t})}\mathbf{x}_{0}$$

Reparametrization (again) Using the reparametrization trick of the forward process, we can reparametrize x_0 as a function of x_t :

$$x_t = \sqrt{\overline{\alpha}_t} x_0 + \sqrt{1 - \overline{\alpha}_t} \epsilon_t \implies x_0 = \frac{1}{\sqrt{\overline{\alpha}_t}} (x_t - \sqrt{1 - \overline{\alpha}_t} \epsilon_t)$$

$$=\frac{\sqrt{\alpha_t}(1-\overline{\alpha}_{t-1})}{(1-\overline{\alpha}_t)}x_t+\frac{\sqrt{\overline{\alpha}_{t-1}}\beta_t}{(1-\overline{\alpha}_t)}x_0$$

By substituing it into the equation, we get:

$$=\frac{\sqrt{\alpha_t}(1-\overline{\alpha}_{t-1})}{(1-\overline{\alpha}_t)}x_t+\frac{\sqrt{\overline{\alpha}_{t-1}}\beta_t}{(1-\overline{\alpha}_t)}\frac{1}{\sqrt{\overline{\alpha}_t}}(x_t-\sqrt{1-\overline{\alpha}_t}\epsilon_t)$$

... don't understand how to solve this part ...

$$\tilde{\mu}(x_t, x_0) = \frac{1}{\sqrt{\alpha_t}} \left(x_t - \frac{1 - \alpha_t}{\sqrt{1 - \overline{\alpha}_t}} \epsilon_t \right)$$

To recap

To recap, we were trying to make sense of the second term (2) in our variational lower bound:

$$L_{VLB} = \mathbb{E}_{q} \left[\underbrace{\sum_{t=2}^{T} D_{KL}(q(x_{t-1}|x_{t},x_{0})||p_{\theta}(x_{t-1}|x_{t}))}_{(2)} - \underbrace{\log p_{\theta}(x_{0}|x_{1})}_{(3)} \right]$$

 $q(x_{t-1}|x_t, x_0)$ is our 'target' gaussian that we want p_{θ} to approximate. We just found a way to calculate it.

$$q(x_{t-1}|x_t, x_0) = \mathcal{N}(x_{t-1}; \tilde{\mu}(x_t, x_0), \tilde{\beta}I)$$

= $\mathcal{N}(x_{t-1}; \frac{1}{\sqrt{\alpha_t}} \left(x_t - \frac{1 - \alpha_t}{\sqrt{1 - \overline{\alpha_t}}} \epsilon_t\right), \tilde{\beta}I)$

To recap, we were trying to make sense of the second term (2) in our variational lower bound:

$$L_{VLB} = \mathbb{E}_{q}[\underbrace{\sum_{t=2}^{T} D_{KL}(q(x_{t-1}|x_{t}, x_{0})||p_{\theta}(x_{t-1}|x_{t}))}_{(2)} - \underbrace{\log p_{\theta}(x_{0}|x_{1})]}_{(3)}]$$

 $p_{\theta}(x_{t-1}|x_t)$ is the distribution we're trying to model.

$$p_{\theta}(x_{t-1}|x_t) = \mathcal{N}(x_{t-1}; \mu_{\theta}(x_t, t), \Sigma_{\theta}(x_t, t))$$
$$= \mathcal{N}(x_{t-1}; \mu_{\theta}(x_t, t), \sigma_t^2 I)$$

In practice (for DDPM), $\sigma_t^2 = \beta_t$ or $\sigma_t^2 = \tilde{\beta}_t$ gives similar results.

N.B.: In a later improved version, [4] proposed to parametrize and **learn the reverse process variance schedule** instead of having it fixed. They proposed an interpolation between β_t and $\tilde{\beta}_t$ by learning a mixing parameter θ :

$$\begin{split} \Sigma_{\theta}(x_t, t) &= \exp(\theta \log \beta_t + (1 - \theta) \log \tilde{\beta}_t) \\ \text{and as we found earlier, } \tilde{\beta}_t &= \frac{1 - \overline{\alpha}_{t-1}}{1 - \overline{\alpha}_t} \beta_t \text{ which means} \\ &= \exp(\theta \log \beta_t + (1 - \theta) \log \left[\frac{1 - \overline{\alpha}_{t-1}}{1 - \overline{\alpha}_t} \beta_t\right]) \end{split}$$

But for now, we'll stick with a reverse process of **fixed variance** schedule.

Training loss

Since the variance of $p_{\theta}(x_{t-1}|x_t)$ is fixed, authors chose to simply minimize the distance between $\mu_{\theta}(x_t, t)$ and $\tilde{\mu}(x_t, t)$:

$$L_t = \mathbb{E}_{\mathbf{x}_0, \epsilon} \left[\frac{1}{2 \left\| \sum_{\theta} (\mathbf{x}_t, t) \right\|_2^2} \left\| \tilde{\mu}_t(\mathbf{x}_t, \mathbf{x}_0) - \mu_{\theta}(\mathbf{x}_t, t) \right\|^2 \right]$$

However, we know that

$$\tilde{\mu}_t(x_t, x_0) = \frac{1}{\sqrt{\alpha_t}} \left(x_t - \frac{1 - \alpha_t}{\sqrt{1 - \overline{\alpha}_t}} \epsilon_t \right)$$

And we also have access to x_t , a_t and $\overline{\alpha}_t$ during the training of μ_{θ} . We can then reparametrize μ_{θ} to only predict the noise ϵ_t instead:

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

The training loss then becomes

$$\begin{split} \mathcal{L}_{t} &= \mathbb{E}_{\mathbf{x}_{0},\epsilon} \left[\frac{1}{2 \left\| \sum_{\theta} \right\|_{2}^{2}} \left\| \tilde{\mu}_{t}(\mathbf{x}_{t}, \mathbf{x}_{0}) - \mu_{\theta}(\mathbf{x}_{t}, t) \right\|^{2} \right] \\ &= \mathbb{E}_{\mathbf{x}_{0},\epsilon} \left[\frac{1}{2 \left\| \sum_{\theta} \right\|_{2}^{2}} \left\| \frac{1}{\sqrt{\alpha_{t}}} (\mathbf{x}_{t} - \frac{1 - \alpha_{t}}{\sqrt{1 - \overline{\alpha}_{t}}} \epsilon_{t}) - \frac{1}{\sqrt{\alpha_{t}}} (\mathbf{x}_{t} - \frac{1 - \alpha_{t}}{\sqrt{1 - \overline{\alpha}_{t}}} \epsilon_{\theta}(\mathbf{x}_{t}, t)) \right\|^{2} \right] \\ &= \mathbb{E}_{\mathbf{x}_{0},\epsilon} \left[\frac{(1 - \alpha_{t})^{2}}{2\alpha_{t}(1 - \overline{\alpha}_{t}) \left\| \sum_{\theta} \right\|_{2}^{2}} \left\| \epsilon_{t} - \epsilon_{\theta}(\mathbf{x}_{t}, t) \right\|^{2} \right] \end{split}$$

Empirically, ignoring the weighting term provides a more stable training for the diffusion model:

$$\begin{split} L_t &= \mathbb{E}_{\mathsf{x}_0,\epsilon} \left[\|\epsilon_t - \epsilon_\theta(\mathsf{x}_t, t)\|^2 \right] \\ &= \mathbb{E}_{\mathsf{x}_0,\epsilon} \left[\left\|\epsilon_t - \epsilon_\theta(\sqrt{\overline{\alpha}_t}\mathsf{x}_0 + \sqrt{1 - \overline{\alpha}_t}\epsilon_t, t)\right\|^2 \right] \end{split}$$
Algorithm 5: Training

repeat

Sample $x_0 \sim q(x_0)$; Sample $t \sim \text{Uniform}(\{1, \dots, T\})$; Sample $\epsilon \sim \mathcal{N}(0, I)$; Take gradient descent step on $\nabla_{\theta} \| \epsilon_t - \epsilon_{\theta} \left(\sqrt{\overline{\alpha}_t} x_0 + \sqrt{1 - \overline{\alpha}_t} \epsilon_t, t \right) \|^2$; until converged;

Algorithm 6: Sampling

```
Sample x_T \sim \mathcal{N}(0, I)

for t = T, \dots, 1 do

if t > 1 then

| Sample z \sim \mathcal{N}(0, I)

else

| z = 0

x_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left( x_t - \frac{1 - \alpha_t}{\sqrt{1 - \bar{\alpha}_t}} \epsilon_{\theta}(x_t, t) \right) + \sigma_t z
```

return x₀

As mentionned before, the overall objective we're trying to minimize is the following:

$$L_{VLB} = \mathbb{E}_{q} \left[\underbrace{\sum_{t=2}^{T} D_{KL}(q(x_{t-1}|x_{t}, x_{0})||p_{\theta}(x_{t-1}|x_{t}))}_{(2)} - \log p_{\theta}(x_{0}|x_{1}) \right]}_{(3)}$$

We know how to minimize (2) with L_t . But, what about (3)? We can't forget about it, otherwise we won't know how to get the image from x_1 to x_0 (the final denoising step).

The final denoising step is modelled by another **independent discrete decoder**, with D being the data dimension (e.g. nb of pixels in an image):

$$p_{ heta}(x_0|x_1) = \prod_{i=1}^{D} \int_{\delta_{-}(x_0^i)}^{\delta_{+}(x_0^i)} \mathcal{N}(x; \mu_{ heta}^i(x_1, 1), \sigma_1^2) dx$$

$$\delta_{-}(x) = \begin{cases} \infty & \text{if } x = 1 \\ x + \frac{1}{255} & \text{if } x < 1 \end{cases} \qquad \delta_{+}(x) = \begin{cases} -\infty & \text{if } x = -1 \\ x - \frac{1}{255} & \text{if } x > -1 \end{cases}$$

In summary:

- 1. Calc. the distribution for the *i*-th pixel of x_0 given the image x_1 .
- 2. Calc. the likelihood of the pixel value x_0^i given the distribution.
- 3. Multiply the likelihood of all pixels to get the final likelihood.



$$L_{VLB} = \mathbb{E}_{q} \left[\sum_{t=2}^{T} D_{KL}(q(x_{t-1}|x_{t}, x_{0})||p_{\theta}(x_{t-1}|x_{t})) - \log p_{\theta}(x_{0}|x_{1}) \right]$$
(2)
(3)

We now have a way to compute (2) and (3), completing the training objective of the foundation of diffusion models:

• (2): Train an independent decoder: $p_{\theta}(x_0|x_1) = \prod_{i=1}^{D} \int_{\delta_{-}(x_0^i)}^{\delta_{+}(x_0^i)} \mathcal{N}(x; \mu_{\theta}^i(x_1, 1), \sigma_1^2) dx$

• (3): Train the diffusion model with the loss:

$$L_t = \mathbb{E}_{x_0,\epsilon} \left[\left\| \epsilon_t - \epsilon_\theta (\sqrt{\overline{\alpha}_t} x_0 + \sqrt{1 - \overline{\alpha}_t} \epsilon_t, t) \right\|^2 \right]$$

We have seen 2 different perspectives on generative modeling using diffusion:

- 1. Score-based generative models: which consists of
 - Modelling the "score" which is the gradients of the log-likelihood of the data ∇_x log p(x) using score matching.
 - 2. Using the gradients with Langevin dynamics to sample new data.
 - 3. Learning from multiple levels of noise perturbation to improve the accuracy of scores in the data space.
 - 4. Learning from an infinite number of noise levels using SDEs.

We have seen 2 different perspectives on generative modeling using diffusion:

- 2. Denoising Diffusion Probabilistic Models (DDPM): which consists of
 - 1. A forward process iteratively adds noise to the data.
 - 2. A reverse process modelled by a model $p_{\theta}(x_{t-1}|x_t)$, that can generate data (images) from a prior distribution.
 - 3. Training the reverse process using a variational lower bound.

Most subsequent developments on diffusion models focus on improving the methods presented above with:

- 1. Faster and more efficient sampling.
- 2. Handling different structures of data.
- 3. More accurate likelihood and density estimation.

Some keywords I haven't covered in this presentation:

- Conditional generation (with/without guidance from classifiers)
- Scale up resolution.
- Latent diffusion models.
- Speed up sampling process.
- Diffusion models for video generation.

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