Accelerating Diffusion Models EECE695D: Efficient ML Systems

Spring 2025

- A generative model, i.e., method to model the probability density $p(\mathbf{x})$
 - "denoising diffusion probabilistic model"

• Like many other generative models, $p(\mathbf{x})$ is generated as a **pushforward** of some easy-to-sample density (e.g., Gaussian)

$$\mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I}_d)$$
 - Gene

• Goal. Learn a good $f(\cdot)$ from samples of **x**

erator $f(\cdot)$ \rightarrow $\mathbf{x} = f(\mathbf{z}) \sim p(\mathbf{x})$

- Challenge.
 - How should we generate the corresponding $\mathbf{z}^{(i)}$ for some sample $\mathbf{x}^{(i)}$? • Also, $f(\cdot)$ is likely to be very complicated

- Idea. There is a straightforward way to model f^{-1} $(\text{use } \mathbf{z}^{(i)} = f^{-1}(\mathbf{x}^{(i)}))$
 - Plus, this f^{-1} can be decomposed into many sub-functions





• Forward diffusion. Adds Gaussian noise gradually

$$q(\mathbf{x}_t | \mathbf{x}_{t-1}) = \mathcal{N}(\mathbf{x}_t | \sqrt{1 - \beta_t} \mathbf{x}_{t-1}, \boldsymbol{\beta}_t \mathbf{I})$$

• Given \mathbf{X}_{0} , we can sample \mathbf{X}_{t} as

$$\mathbf{x}_t = \sqrt{\bar{\alpha}_t} \mathbf{x}_0 + \mathbf{v}$$

 $\mathbf{X}_1 \longrightarrow \mathbf{X}_2 \longrightarrow$



Generative reverse denoising process



X₀

Data

 β_t : analogous to "time" between *t* and t - 1

$$(\cdots)$$
 \longrightarrow $\mathbf{Z} = \mathbf{X}_T$



- Reverse denoising. Want to model q
 - After some math^{*}, one can realize that we can approximate

$$q(\mathbf{x}_{t-1} | \mathbf{x}_t) = \mathcal{N}(\mathbf{x}_{t-1} | \mu_{\theta}(\mathbf{x}_t), \sigma_t^2 \mathbf{I}_d)$$

where the mean $\mu_{\theta}(\mathbf{X}_t)$ can be written as:

$$\mu_{\theta}(\mathbf{x}_{t}, t) = \frac{1}{\sqrt{\alpha_{t}}} \left(\mathbf{x}_{t} - \frac{1 - \alpha_{t}}{\sqrt{1 - \bar{\alpha}_{t}}} \boldsymbol{\epsilon}_{\theta}(\mathbf{x}_{t}, t) \right)$$

squared loss

$$q(\mathbf{x}_{t-1} \,|\, \mathbf{x}_t)$$

Noise model; to be trained from the data

As the model is Gaussian, fitting the distribution is simply training with the



• Training. Train by noise prediction



Algorithm 1 Training

- 1: repeat
- 2: $\mathbf{x}_0 \sim q(\mathbf{x}_0)$
- 3: $t \sim \text{Uniform}(\{1, \ldots, T\})$
- 4: $\boldsymbol{\epsilon} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$
- 5: Take gradient descent step on

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

6: until converged





• Sampling. Step-by-step denoising



Algorithm 2 Sampling

1: $\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ 2: for t = T, ..., 1 do 3: $\mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ if t > 1, else $\mathbf{z} = \mathbf{0}$

4:
$$\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left(\mathbf{x}_t - \frac{1-\alpha_t}{\sqrt{1-\bar{\alpha}_t}} \boldsymbol{\epsilon}_{\theta}(\mathbf{x}_t, t) \right) + \boldsymbol{\sigma}_{\theta}(\mathbf{x}_t, t)$$



 \mathbf{X}_{t-1}

Ho et al., "Denoising Diffusion Probabilistic Models," NeurIPS 2020



Problem

- Goal. We want fast & on-device generation
 - Hopefully video editing as well





Sam Altman 🤣 🕸 @sama · 1h can yall please chill on generating images this is insane our team needs sleep

Q 1,5K 1,1K 🗘 13,1K 1 635K

SAMSUNG









Approaches

- Many different approaches
 - SDE/ODE solvers
 - Reduce the number of denoising steps
 - Deterministic sampler (e.g., DDIM)
 - Distillation
 - Reduce the computational cost of denoising model
 - Compress the model
 - Re-use computed values
 - Parallel sampling

(Presentation 1) (Presentation 2) (Presentation 3)



- Idea. Consider infinitesimal time intervals
 - Recall that the forward diffusion is

$$q(\mathbf{x}_t | \mathbf{x}_{t-1}) = \mathcal{N}(\mathbf{x}_t | \sqrt{1 - \beta_t} \mathbf{x}_{t-1}, \beta_t \mathbf{I})$$

This becomes an stochastic differential equation

$$\begin{aligned} \mathbf{x}_{t} &= \sqrt{1 - \beta(t)\Delta t} \mathbf{x}_{t-1} + \sqrt{\beta(t)\Delta t} \mathcal{N}(0,\mathbf{I}) \\ &\approx \mathbf{x}_{t-1} - \frac{\beta(t)\Delta t}{2} \mathbf{x}_{t-1} + \sqrt{\beta(t)\Delta t} \mathcal{N}(0,\mathbf{I}) \\ &\Rightarrow \quad \mathrm{d}\mathbf{x}_{t} = -\frac{1}{2}\beta(t)\mathbf{x}_{t}\mathrm{d}t + \sqrt{\beta(t)}\mathrm{d}\omega_{t} \end{aligned}$$

ODE solvers





ODE solvers



- The reverse diffusion process can be written as: $d\mathbf{x}_{t} = -\frac{1}{2}\beta(t)\left(\mathbf{x}_{t} + 2\nabla_{\mathbf{x}_{t}}\log q_{t}(\mathbf{x}_{t})\right)dt + \sqrt{\beta(t)}d\bar{\omega}_{t}$
- We can train a neural network which approximates this "score function"
 - Use $q_t(\mathbf{x}_t | \mathbf{x}_0)$ for tractibility
 - The reverse can then be expressed as: $\mathbf{d}\mathbf{x}_t = -\frac{1}{2}\beta(t)\big(\mathbf{x}_t + \mathbf{x}_t^T\big)$
 - Use off-the-shelf SDE solvers.
 - Can also come up with ODE version, which is very fast!

ODE solvers

$$-2\mathbf{s}_{\theta}(\mathbf{x}_{t},t) dt + \sqrt{\beta(t)} d\bar{\omega}_{t}$$

Song et al., "Score-Based Generative Modeling through Stochastic Differential Equations," ICLR 2021



Deterministic sampler

• Idea. We can play with the noise-adding procedure

$$q(\mathbf{x}_t | \mathbf{x}_{t-1}) = \mathcal{N}(\mathbf{x}_t | \sqrt{1 - \beta_t} \mathbf{x}_{t-1}, \beta_t \mathbf{I})$$

- Markov property
- Fixed

- Desired. Want faster "mixing" during forward
 - so that reverting them can be done in fewer steps



Data

(is there a reason why it should be so?)

(can we introduce learnable components?)

Noise

Song et al., "Denoising Diffusion Implicit Models," ICLR 2021



Deterministic sampler

- **DDIM.** No noise-adding during the sampling
 - Theoretical motivations from approximating $q(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{x}_0)$ (not $q(\mathbf{x}_{t-1} | \mathbf{x}_t)$)
 - Same training procedure
 - Faster sampling (1000 \rightarrow 25~40), but slightly weaker image diversity





 \mathbf{X}_{t-1} Song et al., "Denoising Diffusion Implicit Models," ICLR 2021



Distillation

- Idea. Simply distill a multi-step denoiser from a single-step one
 - Luhman & Luhman (2021) distills with the loss

$$L = \frac{1}{2} \mathbb{E}_{\mathbf{x}_T} \| f_{\text{stu}}(\mathbf{x}_T) - f_{\text{tea}}(\mathbf{x}_T) \|_2^2$$

- where $f(\cdot)$ denotes the mean of the estimated Gaussian.
 - f_{tea} is generated by multi-step diffusion



Luhman and Luhman, "Knowledge Distillation in Iterative Generative Models for Improved Sampling Speed," arXiv 2021



Distillation

- Later works find that progressive distillation is beneficial, in general
 - No need to run full number of sampling steps with the original model





Considerations in model compression

- A noteworthy characteristic of diffusion models is their time-dependency
 - The activation distribution changes from timestep to timestep
 - Requires a careful calibration of quantization range



Considerations in feature reuse

- There seems to be much feature redundancy across timesteps
 - Caching and reusing high-level features or attention can save computations at the expense of minimal quality degradation

(a) Examples of Feature Maps

A large teddy bear with a heart is in the garbage A green plate filled with rice and a mixture of sauce on top of it A very ornate, three layered wedding cake in a banquet room



Original

Step20

Step19



Ma et al., "DeepCache: Accelerating Diffusion Models for Free," CVPR 2024

(b) HeatMap for Similarity



Further readings

- Consistency models
 - <u>https://arxiv.org/abs/2303.01469</u>
- Parallel sampling
 - <u>https://arxiv.org/abs/2305.16317</u>
- Early stopping
 - <u>https://arxiv.org/abs/2205.12524</u>

