Deep Generative Learning

Deep Learning for Computer Vision

Valeriya Strizhkova 15.03.23

About myself

Valeriya Strizhkova

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Research topics:

- emotion detection
- computer vision
- deep learning



Outline

- Introduction
 - Overview of deep generative models
 - OpenAl's CLIP
- Variational Autoencoders
- Autoregressive models
- Diffusion Models
- Generative Adversarial Networks
- Evaluation

There are many interesting recent development in deep learning...The most important one, in my opinion, is adversarial training (also called GAN for Generative Adversarial Networks). This, and the variations that are now being proposed, is the most interesting idea in the last 10 years in ML.

Yann LeCun



The Landscape of Deep Generative Learning

Autoregressive Models Normalizing Flows

Variational Autoencoders

Generative Adversarial Networks Energy-based Models Denoising Diffusion Models

pr2022-tutorial-diffusion-models.github.io



A question to ChatGPT

Q: When is it better to use GAN and when diffusion models?

A: The choice between using GANs and diffusion models depends on the specific characteristics of the data being generated and the requirements of the generative modeling task at hand.

In general, GANs are often preferred for tasks that require generating high-quality, visually realistic images, as they have been shown to produce photorealistic images with fine details and textures. GANs are also well-suited for image-to-image translation tasks, such as style transfer or super-resolution, where the goal is to modify an existing image in a particular way while preserving its underlying structure and content.

On the other hand, diffusion models are often preferred for tasks that involve modeling complex distributions or generating samples that exhibit long-range dependencies. Diffusion models have been shown to be effective at modeling complex image and video datasets, and can generate high-quality samples that capture the underlying structure and dynamics of the data. Diffusion models can also be used for tasks such as image denoising or inpainting, where the goal is to remove noise or missing information from an image while preserving its structure and content.

OpenAl's CLIP

- encodes image, and text to similar embeddings
- is trained with contrastive learning, maximizing cosine similarity of corresponding image and text
- CLIP's output image embeddings contain both style and semantics
- is used in DALL-E 1, DALL-E 2, image-text classification



Alec Radford, Jong Wook Kim, Chris Hallacy, Aditya Ramesh, Gabriel Goh, Sandhini Agarwal, Girish Sastry, Amanda Askell, Pamela Mishkin, Jack Clark, Gretchen Krueger, Ilya Sutskever. Learning Transferable Visual Models From Natural Language Supervision. ICML 2021 <u>https://arxiv.org/abs/2103.00020</u> <u>https://openai.com/research/clip</u>

CLIP Architecture

- text and image have separate transformer encoders
- visual encoder is ViT (vision transformer)
- text encoder is GPT-2 transformer
- the fixed-length text embedding is extracted from [EOS] token position
- text token embeddings and image patch embeddings also available
- trained on 256 GPUs for 2 weeks

Alec Radford, Jong Wook Kim, Chris Hallacy, Aditya Ramesh, Gabriel Goh, Sandhini Agarwal, Girish Sastry, Amanda Askell, Pamela Mishkin, Jack Clark, Gretchen Krueger, Ilya Sutskever. Learning Transferable Visual Models From Natural Language Supervision. ICML 2021 <u>https://arxiv.org/abs/2103.00020</u> <u>https://openai.com/research/clip</u>

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- Variational Autoencoders
 - Architecture
 - Maximization of variational lower bound
 - o dVAE
- Autoregressive models
- Diffusion Models
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Autoencoders (non-variational)

Unsupervised method for learning latent features from data without any labels.





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Features need to be **lower dimensional** than the data.





Autoencoders (non-variational)

Unsupervised method for learning latent features from data without any labels.

Features need to be **lower dimensional** than the data.

Limitation: no way to produce any new content

 $L=||\hat{x}-x||_2^2$



• VAE is an autoencoder that learns **latent features** from data and enables **generative process**.

- VAE is an autoencoder that learns **latent features** from data and enables **generative process**.
- Instead of encoding an input as a single point, VAE encodes it as a distribution over the latent space.

Encoder network inputs data x and outputs distribution over latent codes z

Encoder Network

$$q_{\phi}(z|x) = N(\mu_{z|x}, \Sigma_{z|x})$$



Encoder network inputs data x and outputs distribution over latent codes z

Decoder network inputs latent code z and outputs distribution over data x

Encoder Network $q_{\phi}(z|x) = N(\mu_{z|x}, \Sigma_{z|x})$



Decoder Network

$$p_{ heta}(x|z) = N(\mu_{x|z}, \Sigma_{x|z})$$



Decoder neural network represent p(x|z) where x is an image, z is latent features to generate x.

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Decoder neural network represent p(x|z) where x is an image, z is latent features to generate x.

Assume prior p(z) is standard unit diagonal Gaussian.

 $p_{ heta}(x|z) = N(\mu_{x|z}, \Sigma_{x|z})$



Decoder neural network represent p(x|z) where x is an image, z is latent features to generate x.

Assume prior p(z) is standard unit diagonal Gaussian.

How to train this model?

$$p_{ heta}(x|z) = N(\mu_{x|z}, \Sigma_{x|z})$$



Decoder neural network represent p(x|z) where x is an image, z is latent features to generate x.

Assume prior p(z) is standard unit diagonal Gaussian.

How to train this model?

Maximize likelihood of data

$$p_{ heta}(x|z) = N(\mu_{x|z}, \Sigma_{x|z})$$



Maximize likelihood of data

$$p_{ heta}(x) = \int p_{ heta}(x,z) dz = \int p_{ heta}(x|z) p_{ heta}(z) dz$$



Maximize likelihood of data

$$p_{ heta}(x) = \int p_{ heta}(x,z) dz = \int p_{ heta}(x|z) p_{ heta}(z) dz$$

Problem: impossible to integrate over all z



Maximize likelihood of data

Bayes' Rule: $p_{ heta}(x) = rac{p_{ heta}(x|z)p_{ heta}(z)}{p_{ heta}(z|x)}$



Maximize likelihood of data

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Maximize likelihood of data

Bayes' Rule: $p_{ heta}(x) = rac{p_{ heta}(x|z)p_{ heta}(z)}{p_{ heta}(z|x)}$ Problem: no way to compute $p_{ heta}(z|x)$

Solution: train another network (encoder) that learns $q_{\phi}(z|x) pprox p_{ heta}(z|x)$



Encoder network inputs data x and outputs distribution over latent codes z

Decoder network inputs latent code z and outputs distribution over data x

Encoder Network $q_{\phi}(z|x) = N(\mu_{z|x}, \Sigma_{z|x})$



Decoder Network

$$p_{ heta}(x|z) = N(\mu_{x|z}, \Sigma_{x|z})$$



Jointly train encoder q and decoder p

Encoder Network

$$q_{\phi}(z|x) = N(\mu_{z|x}, \Sigma_{z|x})$$



Decoder Network $p_{ heta}(x|z) = N(\mu_{x|z}, \Sigma_{x|z})$



 $\log p_{ heta}(x) = \log rac{p_{ heta}(x|z)p(z)}{p_{ heta}(z|x)}$

$$\log p_{ heta}(x) = \log rac{p_{ heta}(x|z)p(z)}{p_{ heta}(z|x)} = \log rac{p_{ heta}(x|z)p(z)q_{\phi}(z|x)}{p_{ heta}(z|x)q_{\phi}(z|x)}$$

$$egin{aligned} \log p_{ heta}(x) &= \log rac{p_{ heta}(x|z)p(z)}{p_{ heta}(z|x)} = \log rac{p_{ heta}(x|z)p(z)q_{\phi}(z|x)}{p_{ heta}(z|x)q_{\phi}(z|x)} \ &= \log p_{ heta}(x|z) - \log rac{q_{\phi}(z|x)}{p(z)} + \log rac{q_{\phi}(z|x)}{p_{ heta}(z|x)} \end{aligned}$$

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$$egin{aligned} &= \log p_{ heta}(x|z) - \log rac{q_{\phi}(z|x)}{p(z)} + \log rac{q_{\phi}(z|x)}{p_{ heta}(z|x)} \ &= E_{z}[\log p_{ heta}(x|z)] - E_{z}[\log rac{q_{\phi}(z|x)}{p(z)}] + E_{z}[\log rac{q_{\phi}(z|x)}{p_{ heta}(z|x)}] \end{aligned}$$

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$$=E_{z\sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)]-KL(q_{\phi}(z|x),p(z))+KL(q_{\phi}(z|x),p_{ heta}(z|x))$$

Kullback-Leibler Divergence:
$$KL(p,q) = E_{x \sim p} [log rac{p(x)}{q(x)}]$$

$$egin{aligned} &\log p_{ heta}(x) = \log rac{p_{ heta}(x|z)p(z)}{p_{ heta}(z|x)} = \log rac{p_{ heta}(x|z)p(z)q_{\phi}(z|x)}{p_{ heta}(z|x)q_{\phi}(z|x)} \ &= E_{z}[\log p_{ heta}(x|z)] - E_{z}[\log rac{q_{\phi}(z|x)}{p(z)}] + E_{z}[\log rac{q_{\phi}(z|x)}{p_{ heta}(z|x)}] \ &= E_{z\sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)] - KL(q_{\phi}(z|x),p(z)) + KL(q_{\phi}(z|x),p_{ heta}(z|x)) \end{aligned}$$

 $KL \ge 0$ => dropping the last term gives a lower bound on the data likelihood

Kullback-Leibler Divergence:
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$$=E_{z\sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)]-KL(q_{\phi}(z|x),p(z))+KL(q_{\phi}(z|x),p_{ heta}(z|x))$$

$$\log p_{ heta}(x) \geq E_{z \sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)] - KL(q_{\phi}(z|x),p(z))$$

Jointly train **encoder** q and **decoder** p to maximize the **variational lower bound** on the data likelihood

$$\log p_{ heta}(x) \geq E_{z \sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)] - KL(q_{\phi}(z|x),p(z))$$

Encoder Network

$$q_{\phi}(z|x) = N(\mu_{z|x}, \Sigma_{z|x})$$

Decoder Network

$$p_{ heta}(x|z) = N(\mu_{x|z}, \Sigma_{x|z})$$




Closed form solution when q_{ϕ} is diagonal Gaussian and p is unit Gaussian:

$$egin{aligned} &\log p_{ heta}(x) \geq E_{z \sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)] - KL(q_{\phi}(z|x),p(z)) \ &- KL(q_{\phi}(z|x),p(z)) = \int_{Z} q_{\phi}(z|x) \log rac{p(z)}{q_{\phi}(z|x)} dz \ &= \int_{Z} N(z,\mu_{z|x},\Sigma_{z|x}) \log rac{N(z,0,I)}{N(z,\mu_{z|x},\Sigma_{z|x})} dz \ &= \sum_{j=1}^{J} (1 + \log((\sum_{z|x})_{j}^{2}) - (\mu_{z|x})_{j}^{2} - (\Sigma_{z|x})_{j}^{2}) \end{aligned}$$

Closed form solution when q_{ϕ} is diagonal Gaussian and p is unit Gaussian:

$$\log p_{ heta}(x) \geq E_{z \sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)] - KL(q_{\phi}(z|x),p(z))$$

 $E_{z\sim q_{\phi}(z|x)}[\log p_{ heta}(x|z)]$ is data reconstruction term

Train by maximize the variational lower bound.

 $E_{z\sim q_{\phi}(z|x)}[\log p_{\Theta}(x|z)] - KL(q_{\phi}(z|x),p(z))$

1. The input is **encoded** as distribution over the latent space



Train by maximize the variational lower bound.

 $E_{z\sim q_{\phi}(z|x)}[\log p_{\Theta}(x|z)] - KL(q_{\phi}(z|x),p(z))$

- 1. The input is **encoded** as distribution over the latent space
- 2. Encoder output should match prior p(z)



Train by maximize the variational lower bound.

 $E_{z\sim q_{\phi}(z|x)}[\log p_{\Theta}(x|z)] - KL(q_{\phi}(z|x),p(z))$

- 1. The input is **encoded** as distribution over the latent space
- 2. Encoder output should match prior p(z)
- 3. A point from the latent space is sampled from that distribution



Train by maximize the variational lower bound.

 $E_{z\sim q_{\phi}(z|x)}[\log p_{\Theta}(x|z)] - KL(q_{\phi}(z|x),p(z))$

- 1. The input is **encoded** as distribution over the latent space
- 2. Encoder output should match prior p(z)
- 3. A point from the latent space is sampled from that distribution
- 4. The sampled point is **decoded**



Sample from z



Train by maximize the variational lower bound.

 $E_{z\sim q_{\phi}(z|x)}[\log p_{\Theta}(x|z)] - KL(q_{\phi}(z|x),p(z))$

- 1. The input is **encoded** as distribution over the latent space
- 2. Encoder output should match prior p(z)
- 3. A point from the latent space is sampled from that distribution
- 4. The sampled point is **decoded**
- 5. The reconstruction error is computed



Sample from z



Learned data manifold for generative models with two-dimensional latent space:



В n n Б з з

Discrete Variational Auto-Encoder (dVAE)

- introduced in VQ-VAE 1 [1] and VQ-VAE-2 [2]
- image encoder maps to latent 32x32 grid of embeddings
- vector quantization maps to 8k code words
- decoder maps from quantized grid to the image
- copy gradients from decoder input z to the encoder output



[1] Aaron van den Oord, Oriol Vinyals, Koray Kavukcuoglu. Neural Discrete Representation Learning. NeurIPS 2017 https://arxiv.org/abs/1711.00937

[2] Ali Razavi, Aaron van den Oord, Oriol Vinyals. Generating Diverse High-Fidelity Images with VQ-VAE-2. NeurIPS 2019 https://arxiv.org/abs/1906.00446

Outline

- Introduction
- Variational Autoencoders
- Autoregressive models
 - Image generation
 - PixelCNN
 - Text-to-image generation
 - DALL-E 1
 - Paella
 - Text-to-video generation
 - Phenaki
- Diffusion Models
- Generative Adversarial Networks
- Evaluation

Autoregressive Models

Output • • • • • • • • • • • • • • • • • •

Hidden Layer

Layer OOOOOOOOOOOOOOOOOOOO

Pixel Recurrent Neural Networks





DALL-E 1

- introduced by OpenAl
- generates 256×256 images from text via dVAE inspired by VQ-VAE-2
- autoregressive-ly generates image tokens from textual tokens on a discrete latent space

DALL-E 1 Training

- train encoder and decoder image of image into 32x32 grid of 8k possible code word tokens (dVAE)
- 2. concatenate encoded text tokens with image tokens into single array
- 3. train to predict next image token from the preceding tokens (autoregressive transformer)
- 4. discard the image encoder, keep only image decoder and next token predictor

DALL-E 1 Prediction

- 1. encode input text to text tokens
- 2. iteratively predict next image token from the learned codebook
- 3. decode the image tokens using dVAE decoder
- 4. select the best image using CLIP model ranker

DALL-E 1 Discrete Variational Auto-Encoder (dVAE)

- instead of copying gradients annealing (categorical reparameterization with gumbel-softmax)
- promote codebook utilization using higher KL-divergence weight
- decoder is conv2d, decoder block (4x relu + conv), upsample (tile bigger array), repeat

DALL-E 1 Results



accordion.

sweater walking a dog

(a) a tapir made of accordion. (b) an illustration of a baby (c) a neon sign that reads (d) the exact same cat on the a tapir with the texture of an hedgehog in a christmas "backprop". a neon sign that top as a sketch on the bottom reads "backprop". backprop neon sign

Aditya Ramesh, Mikhail Pavlov, Gabriel Goh, Scott Gray, Chelsea Voss, Alec Radford, Mark Chen, Ilya Sutskever. Zero-Shot Text-to-Image Generation. PMLR 2021 https://arxiv.org/abs/2102.12092

DALL-E 1 Results



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Paella



An oil painting of a cute german shepard.



A crayon drawing of barack obama.



A cat dressed as a medieval knight.



ship floating on the water.

High quality photograph of the milkyway at Zugspitze.



Image Variations



Latent Space Interpolations



A beautiful bavarian village with the alps in the backtground.







Outpainting

Dominic Rampas, Pablo Pernias, Elea Zhong, Marc Aubreville. Fast Text-Conditional Discrete Denoising on Vector-Quantized Latent Spaces. 2022 https://arxiv.org/abs/2211.07292



Dominic Rampas, Pablo Pernias, Elea Zhong, Marc Aubreville. Fast Text-Conditional Discrete Denoising on Vector-Quantized Latent Spaces. 2022 https://arxiv.org/abs/2211.07292



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Paella



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Phenaki

The water is magical



Prompts used:

A photorealistic teddy bear is swimming in the ocean at San Francisco The teddy bear goes under water The teddy bear keeps swimming under the water with colorful fishes A panda bear is swimming under water

Chilling on the beach



Prompts used:

- A teddy bear diving in the ocean A teddy bear emerges from the water A teddy bear walks on the beach
- Camera zooms out to the teddy bear in the campfire by the beach

Fireworks on the spacewalk



Prompts used: Side view of an astronaut is walking through a puddle on mars The astronaut is dancing on mars The astronaut walks his dog on mars The astronaut and his dog watch fireworks

Ruben Villegas, Mohammad Babaeizadeh, Pieter-Jan Kindermans, Hernan Moraldo, Han Zhang, Mohammad Taghi Saffar, Santiago Castro, Julius Kunze, Dumitru Erhan. Phenaki: Variable Length Video Generation From Open Domain Textual Description. 2022 <u>https://arxiv.org/abs/2210.02399</u> <u>https://phenaki.video</u>

Phenaki



Ruben Villegas, Mohammad Babaeizadeh, Pieter-Jan Kindermans, Hernan Moraldo, Han Zhang, Mohammad Taghi Saffar, Santiago Castro, Julius Kunze, Dumitru Erhan. Phenaki: Variable Length Video Generation From Open Domain Textual Description. 2022 <u>https://arxiv.org/abs/2210.02399</u> <u>https://phenaki.video</u>

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 - Denoising diffusion models
 - Text-to-image generation
 - DALL-E 2
 - Stable Diffusion
 - Text-to-video generation
 - Make-A-Video
- Generative Adversarial Networks
- Evaluation



Diffusion Models

Diffusion models are inspired by non-equilibrium thermodynamics. They define a Markov chain of diffusion steps to slowly add random noise to data and then learn to reverse the diffusion process to construct desired data samples from the noise. Unlike VAE or flow models, diffusion models are learned with a fixed procedure and the latent variable has high dimensionality (same as the original data).

Slide Credits

CVPR 2022 Tutorial

Denoising Diffusion-based Generative Modeling

by Karsten Kreis, Ruiqi Gao and Arash Vahdat

https://cvpr2022-tutorial-diffusion-models.github.io

Denoising Diffusion Models

Learning to generate by denoising

Denoising diffusion models consist of two processes:

- Forward diffusion process that gradually adds noise to input
- Reverse denoising process that learns to generate data by denoising

Forward diffusion process (fixed)

Data



Noise

Reverse denoising process (generative)

Sohl-Dickstein et al., Deep Unsupervised Learning using Nonequilibrium Thermodynamics, ICML 2015 Ho et al., Denoising Diffusion Probabilistic Models, NeurIPS 2020 Song et al., Score-Based Generative Modeling through Stochastic Differential Equations, ICLR 2021

Forward Diffusion Process

The formal definition of the forward process in T steps:



Diffusion Kernel



 β_t values schedule (i.e., the noise schedule) is designed such that $\bar{\alpha}_T \to 0$ and $q(\mathbf{x}_T | \mathbf{x}_0) \approx \mathcal{N}(\mathbf{x}_T; \mathbf{0}, \mathbf{I})$

What happens to a distribution in the forward diffusion?

So far, we discussed the diffusion kernel $q(\mathbf{x}_t|\mathbf{x}_0)$ but what about $q(\mathbf{x}_t)$?



We can sample $\mathbf{x}_t \sim q(\mathbf{x}_t)$ by first sampling $\mathbf{x}_0 \sim q(\mathbf{x}_0)$ and then sampling $\mathbf{x}_t \sim q(\mathbf{x}_t | \mathbf{x}_0)$ (i.e., ancestral sampling).

Generative Learning by Denoising

Recall, that the diffusion parameters are designed such that $q(\mathbf{x}_T) \approx \mathcal{N}(\mathbf{x}_T; \mathbf{0}, \mathbf{I})$





In general, $q(\mathbf{x}_{t-1}|\mathbf{x}_t) \propto q(\mathbf{x}_{t-1})q(\mathbf{x}_t|\mathbf{x}_{t-1})$ is intractable.

Can we approximate $q(\mathbf{x}_{t-1}|\mathbf{x}_t)$? Yes, we can use a Normal distribution if β_t is small in each forward diffusion step.

Reverse Denoising Process

Formal definition of forward and reverse processes in T steps:

Data



Learning Denoising Model Variational upper bound

For training, we can form variational upper bound that is commonly used for training variational autoencoders:

$$\mathbb{E}_{q(\mathbf{x}_0)}\left[-\log p_{\theta}(\mathbf{x}_0)\right] \le \mathbb{E}_{q(\mathbf{x}_0)q(\mathbf{x}_{1:T}|\mathbf{x}_0)}\left[-\log \frac{p_{\theta}(\mathbf{x}_{0:T})}{q(\mathbf{x}_{1:T}|\mathbf{x}_0)}\right] =: L$$

Sohl-Dickstein et al. ICML 2015 and Ho et al. NeurIPS 2020 show that:

$$L = \mathbb{E}_q \left[\underbrace{D_{\mathrm{KL}}(q(\mathbf{x}_T | \mathbf{x}_0) | | p(\mathbf{x}_T))}_{L_T} + \sum_{t>1} \underbrace{D_{\mathrm{KL}}(q(\mathbf{x}_{t-1} | \mathbf{x}_t, \mathbf{x}_0) | | p_{\theta}(\mathbf{x}_{t-1} | \mathbf{x}_t))}_{L_{t-1}} \underbrace{-\log p_{\theta}(\mathbf{x}_0 | \mathbf{x}_1))}_{L_0} \right]$$

where $q(\mathbf{x}_{t-1}|\mathbf{x}_t,\mathbf{x}_0)$ is the tractable posterior distribution:

$$q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0) = \mathcal{N}(\mathbf{x}_{t-1}; \tilde{\mu}_t(\mathbf{x}_t, \mathbf{x}_0), \tilde{\beta}_t \mathbf{I}),$$

where $\tilde{\mu}_t(\mathbf{x}_t, \mathbf{x}_0) := \frac{\sqrt{\bar{\alpha}_{t-1}}\beta_t}{1 - \bar{\alpha}_t} \mathbf{x}_0 + \frac{\sqrt{1 - \beta_t}(1 - \bar{\alpha}_{t-1})}{1 - \bar{\alpha}_t} \mathbf{x}_t$ and $\tilde{\beta}_t := \frac{1 - \bar{\alpha}_{t-1}}{1 - \bar{\alpha}_t} \beta_t$
Parameterizing the Denoising Model

Since both $q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0)$ and $p_{\theta}(\mathbf{x}_{t-1}|\mathbf{x}_t)$ are Normal distributions, the KL divergence has a simple form:

$$L_{t-1} = D_{\mathrm{KL}}(q(\mathbf{x}_{t-1}|\mathbf{x}_t, \mathbf{x}_0)||p_{\theta}(\mathbf{x}_{t-1}|\mathbf{x}_t)) = \mathbb{E}_q\left[\frac{1}{2\sigma_t^2}||\tilde{\mu}_t(\mathbf{x}_t, \mathbf{x}_0) - \mu_{\theta}(\mathbf{x}_t, t)||^2\right] + C$$

Recall that $\mathbf{x}_t = \sqrt{\bar{lpha}_t} \ \mathbf{x}_0 + \sqrt{(1-\bar{lpha}_t)} \ \epsilon$. Ho et al. NeurIPS 2020 observe that:

$$\tilde{\mu}_t(\mathbf{x}_t, \mathbf{x}_0) = \frac{1}{\sqrt{1 - \beta_t}} \left(\mathbf{x}_t - \frac{\beta_t}{\sqrt{1 - \bar{\alpha}_t}} \epsilon \right)$$

They propose to represent the mean of the denoising model using a *noise-prediction* network:

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

With this parameterization

$$L_{t-1} = \mathbb{E}_{\mathbf{x}_0 \sim q(\mathbf{x}_0), \epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[\frac{\beta_t^2}{2\sigma_t^2 (1 - \beta_t)(1 - \bar{\alpha}_t)} ||\epsilon - \epsilon_\theta (\underbrace{\sqrt{\bar{\alpha}_t} \ \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \ \epsilon}_{\mathbf{X}_t} \epsilon, t) ||^2 \right] + C$$

Training Objective Weighting Trading likelihood for perceptual quality

$$L_{t-1} = \mathbb{E}_{\mathbf{x}_0 \sim q(\mathbf{x}_0), \epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \left[\underbrace{\frac{\beta_t^2}{2\sigma_t^2 (1 - \beta_t)(1 - \bar{\alpha}_t)}}_{\lambda_t} ||\epsilon - \epsilon_\theta (\sqrt{\bar{\alpha}_t} \ \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \ \epsilon, t)||^2 \right]$$

The time dependent λ_t ensures that the training objective is weighted properly for the maximum data likelihood training. However, this weight is often very large for small t's.

<u>Ho et al. NeurIPS 2020</u> observe that simply setting $\lambda_t = 1$ improves sample quality. So, they propose to use:

$$L_{\text{simple}} = \mathbb{E}_{\mathbf{x}_0 \sim q(\mathbf{x}_0), \epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I}), t \sim \mathcal{U}(1, T)} \left[||\epsilon - \epsilon_{\theta} (\underbrace{\sqrt{\bar{\alpha}_t} \ \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \ \epsilon}_{\mathbf{x}_t} \epsilon, t)||^2 \right]$$

For more advanced weighting see Choi et al., Perception Prioritized Training of Diffusion Models, CVPR 2022.

Summary Training and Sample Generation

Algorithm 1 Training	Algorithm 2 Sampling
1: repeat 2: $\mathbf{x}_0 \sim q(\mathbf{x}_0)$ 3: $t \sim \text{Uniform}(\{1, \dots, T\})$ 4: $\boldsymbol{\epsilon} \sim \mathcal{N}(0, \mathbf{I})$ 5: Take gradient descent step on $\nabla_{\theta} \ \boldsymbol{\epsilon} - \boldsymbol{\epsilon}_{\theta} (\sqrt{\overline{\alpha}_t} \mathbf{x}_0 + \sqrt{1 - \overline{\alpha}_t} \boldsymbol{\epsilon}, t) \ ^2$ 6: until converged	1: $\mathbf{x}_T \sim \mathcal{N}(0, \mathbf{I})$ 2: for $t = T,, 1$ do 3: $\mathbf{z} \sim \mathcal{N}(0, \mathbf{I})$ 4: $\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left(\mathbf{x}_t - \frac{1-\alpha_t}{\sqrt{1-\overline{\alpha}_t}} \boldsymbol{\epsilon}_{\theta}(\mathbf{x}_t, t) \right) + \sigma_t \mathbf{z}$ 5: end for 6: return \mathbf{x}_0

Implementation Considerations Network Architectures

Diffusion models often use U-Net architectures with ResNet blocks and self-attention layers to represent $\epsilon_{\theta}(\mathbf{x}_t, t)$



Time representation: sinusoidal positional embeddings or random Fourier features.

Time features are fed to the residual blocks using either simple spatial addition or using adaptive group normalization layers. (see <u>Dharivwal and Nichol NeurIPS 2021</u>)

Summary Denoising Diffusion Probabilistic Models

In this part, we reviewed denoising diffusion probabilistic models.

The model is trained by sampling from the forward diffusion process and training a denoising model to predict the noise.

We discussed how the forward process perturbs the data distribution or data samples.

The devil is in the details:

- Network architectures
- Objective weighting
- Diffusion parameters (i.e., noise schedule)

See <u>"Elucidating the Design Space of Diffusion-Based Generative Models" by Karras et al.</u> for important design decisions.

DALL-E 2

- introduced by OpenAI
- generates 1024 x 1024 images from text using diffusion models.



vibrant portrait painting of Salvador Dalí with a robotic half face

a shiba inu wearing a beret and black turtleneck

a close up of a handpalm with leaves growing from it

Aditya Ramesh, Prafulla Dhariwal, Alex Nichol, Casey Chu, Mark Chen.

Hierarchical Text-Conditional Image Generation with CLIP Latents. 2022 https://arxiv.org/abs/2204.06125

DALL-E 2

- 1. generates a CLIP model text embedding for text caption
- 2. "prior" network generates CLIP image embedding from text embedding
- 3. diffusion decoder generates image from the image embedding



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DALL-E 2 training

- 1. generates a CLIP model text embedding for text caption
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DALL-E 2 training

- Can vary images while preserving style and semantics in the embeddings
- Authors found diffusion models more efficient and higher quality compared to autoregressive



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Hierarchical Text-Conditional Image Generation with CLIP Latents. 2022 https://arxiv.org/abs/2204.06125



Figure 4: Variations between two images by interpolating their CLIP image embedding and then decoding with a diffusion model. We fix the decoder seed across each row. The intermediate variations naturally blend the content and style from both input images.

DALL-E 2 evaluation results



Aditya Ramesh, Prafulla Dhariwal, Alex Nichol, Casey Chu, Mark Chen.

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DALL-E 2 Limitations



A family dining for Christmas

A hyperrealistic man's face smiling with a pair of sunglasses A green-eyed woman with freckles showing happines

A green-eyed woman with freckles showing disgust

Midjourney



Capture the essence of a young footballer wearing the Paris Saint-Germain cinematic uniform through an extreme close-up portrait, focusing on the intricate details of her face and expressions, very detailed, Octane Render, cinematic, digital art



animated blonde blue eyes love tattoo soft male guy enhanced by ai smile defined --q 2 --s 750

surprised and happy man realistic

Midjourney



eiffel tower in space with an alien on top



model and robot in fashion show, public, cyberpunk,small robot ,red road , cyberpunk,8k rendering

Latent Diffusion Models



Text-to-Image Synthesis on LAION. 1.45B Model.



Robin Rombach, Andreas Blattmann, Dominik Lorenz, Patrick Esser, Björn Ommer.

High-Resolution Image Synthesis with Latent Diffusion Models. 2022 https://arxiv.org/abs/2112.10752

Latent Diffusion Models



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Latent Diffusion Models



Robin Rombach, Andreas Blattmann, Dominik Lorenz, Patrick Esser, Björn Ommer.

High-Resolution Image Synthesis with Latent Diffusion Models. 2022 https://arxiv.org/abs/2112.10752

Stable Diffusion

Prompt: realistic smiling woman



Make-A-Video









A teddy bear painting a portrait

Robot dancing in times square

Cat watching TV with a remote in hand

A fluffy baby sloth with an orange knitted hat trying to figure out a laptop close up highly detailed studio lighting screen reflecting in its eye

Uriel Singer, Adam Polyak, Thomas Hayes, Xi Yin, Jie An, Songyang Zhang, Qiyuan Hu, Harry Yang, Oron Ashual, Oran Gafni, Devi Parikh, Sonal Gupta, Yaniv Taigman.

Make-A-Video: Text-to-Video Generation without Text-Video Data. 2022 https://arxiv.org/abs/2209.14792

https://makeavideo.studio

Make-A-Video architecture



Figure 2: Make-A-Video high-level architecture. Given input text x translated by the prior P into an image embedding, and a desired frame rate fps, the decoder D^t generates 16 64 × 64 frames, which are then interpolated to a higher frame rate by \uparrow_F , and increased in resolution to 256 × 256 by SR_l^t and 768 × 768 by SR_h , resulting in a high-spatiotemporal-resolution generated video \hat{y} .

Uriel Singer, Adam Polyak, Thomas Hayes, Xi Yin, Jie An, Songyang Zhang, Qiyuan Hu, Harry Yang, Oron Ashual, Oran Gafni, Devi Parikh, Sonal Gupta, Yaniv Taigman.

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Outline

- Introduction
- Variational Autoencoders
- Autoregressive models
- Diffusion Models
- Generative Adversarial Networks
 - Architecture and training objective
 - Image generation
 - Video generation
 - Problems of GANs
- Evaluation

• Setup: Assume we have data x_i drawn from distribution $p_{data}(x)$. Want to sample from p_{data} .

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Jointly train generator G and discriminator D with a minimax game

$$\min \max_{\mathbf{G}} \left[E_{x \sim p_{data}} \left[\log \mathbf{D}(x) \right] + E_{\mathbf{Z} \sim p(\mathbf{Z})} \left[\log(1 - \mathbf{D}(\mathbf{G}(\mathbf{Z}))) \right]
ight)$$



Jointly train generator G and discriminator D with a minimax game





Jointly train generator G and discriminator D with a minimax game

$$\min_{\mathbf{G}} \max_{\mathbf{D}} (E_{x \sim p_{data}} [\log \mathbf{D}(x)] + E_{\mathbf{Z} \sim p(\mathbf{Z})} [\log(1 - \mathbf{D}(\mathbf{G}(\mathbf{Z})))])$$



Jointly train generator G and discriminator D with a minimax game



Jointly train generator G and discriminator D with a minimax game

$$\begin{split} \min \max & (E_{x \sim p_{data}} \left[\log \mathbf{D}(x) \right] + E_{\mathbf{Z} \sim p(\mathbf{Z})} \left[\log(1 - \mathbf{D}(\mathbf{G}(\mathbf{Z}))) \right]) \\ & = \min \max \mathbf{V}(\mathbf{G}, \mathbf{D}) \\ & \mathbf{G} \quad \mathbf{D} \end{split}$$

Train G and D using alternating gradient updates:

1. Update
$$\mathbf{D} = \mathbf{D} + \alpha_{\mathbf{D}} \frac{\delta \mathbf{V}}{\delta \mathbf{D}}$$

2. Update $\mathbf{G} = \mathbf{G} + \alpha_{\mathbf{G}} \frac{\delta \mathbf{V}}{\delta \mathbf{G}}$

Generative Adversarial Networks: vanishing gradient

$$\min \max_{G} V(G,D) = \min_{G} \max_{D} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight])$$



Generative Adversarial Networks: vanishing gradient

$$\min \max_{G} V(G,D) = \min_{G} \max_{D} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight])$$

$$\begin{aligned} \nabla_{\Theta_G} V(G,D) &= \nabla_{\Theta_G} E_{z \sim q(z)} \left[log(1 - D(G(z))) \right] \\ \nabla_a \log(1 - \sigma(a)) &= \frac{-\nabla_a \sigma(a)}{1 - \sigma(a)} = \frac{-\sigma(a)(1 - \sigma(a))}{1 - \sigma(a)} = -\sigma(a) = -D(G(z)) \end{aligned}$$

$$& \text{Gradient goes to 0 if D is confident, i.e. } D(G(z)) \rightarrow 0 \\ & \text{Minimize } \left[-E_{z \sim p(z)} \left[\log(D(G(z))) \right] \right] \text{ for generator } \end{aligned}$$

Ian Goodfellow, Jean Pouget-Abadie, Mehdi Mirza, Bing Xu, David Warde-Farley, Sherjil Ozair, Aaron Courville, Yoshua Bengio. Generative Adversarial Nets. Advances in Neural Information Processing Systems (NeurIPS) 2014.

0.0

0.2

0.4

D(G(z))/D(x)

0.6

0.8

Generative Adversarial Networks: Optimality

 $\min_{G} \max_{D} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight])$

Generative Adversarial Networks: Optimality

$$egin{aligned} &\min \max_{G} (E_{x \sim p_{data}} [\log D(x)] + E_{z \sim p(z)} [\log(1 - D(G(z)))]) \ &= \min \max_{G} (E_{x \sim p_{data}} [\log D(x)] + E_{x \sim p_G} [\log(1 - D(x))]) \ &G \quad D \end{aligned}$$

Generative Adversarial Networks: Optimality

$$egin{aligned} &\min \max_G (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight] \ &= \min \max_G (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{x \sim p_G} \left[\log(1 - D(x))
ight]) \ &= \min \max_G \int_X (p_{data}(x) \log D(x) + p_G(x) \log(1 - D(x))) dx \end{aligned}$$
$$egin{aligned} &\min \max_{G} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight] \ &= \min \max_{G} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{x \sim p_G} \left[\log(1 - D(x))
ight]) \ &= \min \max_{G} \int_X (p_{data}(x) \log D(x) + p_G(x) \log(1 - D(x))) dx \ &= \min_{G} \int_X \max_D (p_{data}(x) \log D(x) + p_G(x) \log(1 - D(x))) dx \end{aligned}$$

$$egin{aligned} &\min \max_{G} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight]) \ &= \min \max_{G} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{x \sim p_G} \left[\log(1 - D(x))
ight]) \ &= \min \max_{G} \int_X (p_{data}(x) \log D(x) + p_G(x) \log(1 - D(x))) dx \ &= \min_{G} \int_X \max_D (p_{data}(x) \log D(x) + p_G(x) \log(1 - D(x))) dx \end{aligned}$$

Ian Goodfellow, Jean Pouget-Abadie, Mehdi Mirza, Bing Xu, David Warde-Farley, Sherjil Ozair, Aaron Courville, Yoshua Bengio. Generative Adversarial Nets. Advances in Neural Information Processing Systems (NeurIPS) 2014.

 $f(y) = a \log y + b \log(1-y)$

$$egin{aligned} &\min\max_G (E_{x\sim p_{data}}\left[\log D(x)
ight] + E_{z\sim p(z)}\left[\log(1-D(G(z)))
ight]) \ &=\min\max_G (E_{x\sim p_{data}}\left[\log D(x)
ight] + E_{x\sim p_G}\left[\log(1-D(x))
ight]) \ &=\min\max_G \int_X (p_{data}(x)\log D(x) + p_G(x)\log(1-D(x)))dx \ &=\min_G \int_X \max_D (p_{data}(x)\log D(x) + p_G(x)\log(1-D(x)))dx \end{aligned}$$

$$egin{aligned} f(y) &= a\log y + b\log(1-y) \ f'(y) &= rac{a}{y} - rac{b}{1-y} \end{aligned}$$

$$egin{aligned} &\min\max_G (E_{x\sim p_{data}}\left[\log D(x)
ight]+E_{z\sim p(z)}\left[\log(1-D(G(z)))
ight])\ &=\min\max_G (E_{x\sim p_{data}}\left[\log D(x)
ight]+E_{x\sim p_G}\left[\log(1-D(x))
ight])\ &G\quad D\ &=\min\max_G \int_X (p_{data}(x)\log D(x)+p_G(x)\log(1-D(x)))dx\ &=\min_G \int_X \max_D (p_{data}(x)\log D(x)+p_G(x)\log(1-D(x)))dx \end{aligned}$$

$$egin{aligned} f(y) &= a\log y + b\log(1-y) \ f'(y) &= rac{a}{y} - rac{b}{1-y} \ f'(y) &= 0 \Leftrightarrow y = rac{a}{a+b} \end{aligned}$$

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ight])\ &=\min\max_G (E_{x\sim p_{data}}\left[\log D(x)
ight]+E_{x\sim p_G}\left[\log(1-D(x))
ight])\ &G\ D\ &=\min\max_G \int_X (p_{data}(x)\log D(x)+p_G(x)\log(1-D(x)))dx\ &=\min_G \int_X \max_D (p_{data}(x)\log D(x)+p_G(x)\log(1-D(x)))dx \end{aligned}$$

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Optimal Discriminator:

$$D^*_G(x) = rac{p_{data}(x)}{p_{data}(x) + p_G(x)}$$

$$egin{aligned} &\min &\max (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log (1 - D(G(z)))
ight]) \ &= &\min &\max (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{x \sim p_G} \left[\log (1 - D(x))
ight]) \ &G &D \ &= &\min &\int_X (p_{data}(x) \log D^*_G(x) + p_G(x) \log (1 - D^*_G(x))) dx \ &G \end{aligned}$$

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ight]) \ &= \min \max_G (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{x \sim p_G} \left[\log(1 - D(x))
ight]) \ &= \min_G \int_X (p_{data}(x) \log D^*_G(x) + p_G(x) \log(1 - D^*_G(x))) dx \ &= \min_G \int_X (p_{data}(x) \log rac{p_{data}(x)}{p_{data}(x) + p_G(x)} + p_G(x) \log rac{p_G(x)}{p_{data}(x) + p_G(x)}) dx \end{aligned}$$

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ight]) \ &= \min \int_X (p_{data}(x) \log D^*_G(x) + p_G(x) \log(1 - D^*_G(x))) dx \ &G \ &= \min_G \int_X (p_{data}(x) \log rac{p_{data}(x)}{p_{data}(x) + p_G(x)} + p_G(x) \log rac{p_G(x)}{p_{data}(x) + p_G(x)}) dx \ &= \min_G (E_{x \sim p_{data}} \left[\log rac{p_{data}(x)}{p_{data}(x) + p_G(x)}
ight] + E_{x \sim p_G} \left[\log rac{p_G(x)}{p_{data}(x) + p_G(x)}
ight]) \end{aligned}$$

$$egin{aligned} &\min_{G} \max_{D} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight]) \ &= \min_{G} \int_{X} (p_{data}(x) \log D_{G}^{*}(x) + p_{G}(x) \log(1 - D_{G}^{*}(x))) dx \ &= \min_{G} \int_{X} (p_{data}(x) \log rac{p_{data}(x)}{p_{data}(x) + p_{G}(x)} + p_{G}(x) \log rac{p_{G}(x)}{p_{data}(x) + p_{G}(x)}) dx \ &= \min_{G} (E_{x \sim p_{data}} \left[\log rac{p_{data}(x)}{p_{data}(x) + p_{G}(x)}
ight] + E_{x \sim p_{G}} \left[\log rac{p_{G}(x)}{p_{data}(x) + p_{G}(x)}
ight]) \ &= \min_{G} (E_{x \sim p_{data}} \left[\log rac{2 imes p_{data}(x)}{p_{data}(x) + p_{G}(x)}
ight] + E_{x \sim p_{G}} \left[\log rac{2 imes p_{G}(x)}{p_{data}(x) + p_{G}(x)}
ight] - \log 4) \end{aligned}$$

$$egin{aligned} &\min_{G} \max_{D}(E_{x \sim p_{data}}[\log D(x)] + E_{z \sim p(z)}[\log(1 - D(G(z)))]) \ &= \min_{G}(E_{x \sim p_{data}}[\lograc{2 imes p_{data}(x)}{p_{data}(x) + p_{G}(x)}] + E_{x \sim p_{G}}[\lograc{2 imes p_{G}(x)}{p_{data}(x) + p_{G}(x)}] - \log 4) \end{aligned}$$

$$egin{aligned} &\min_{G} \max_{D}(E_{x \sim p_{data}}[\log D(x)] + E_{z \sim p(z)}[\log(1 - D(G(z)))]) \ &= \min_{G}(E_{x \sim p_{data}}[\lograc{2 imes p_{data}(x)}{p_{data}(x) + p_{G}(x)}] + E_{x \sim p_{G}}[\lograc{2 imes p_{G}(x)}{p_{data}(x) + p_{G}(x)}] - \log 4) \end{aligned}$$

Kullback-Leibler Divergence:
$$\,KL(p,q)=E_{x\sim p}[lograc{p(x)}{q(x)}]$$

$$egin{aligned} &\min \max_G (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log (1 - D(G(z)))
ight]) \ &= \min_G (E_{x \sim p_{data}} \left[\log rac{2 imes p_{data}(x)}{p_{data}(x) + p_G(x)}
ight] + E_{x \sim p_G} \left[\log rac{2 imes p_G(x)}{p_{data}(x) + p_G(x)}
ight] - \log 4) \ &= \min_G (KL(p_{data}, rac{p_{data} + p_G}{2}) + KL(p_G, rac{p_{data} + p_G}{2}) - \log 4) \end{aligned}$$

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ight]) \ &= \min_G (E_{x \sim p_{data}} \left[\log rac{2 imes p_{data}(x)}{p_{data}(x) + p_G(x)}
ight] + E_{x \sim p_G} \left[\log rac{2 imes p_G(x)}{p_{data}(x) + p_G(x)}
ight] - \log 4) \ &= \min_G (KL(p_{data}, rac{p_{data} + p_G}{2}) + KL(p_G, rac{p_{data} + p_G}{2}) - \log 4) \end{aligned}$$

Jensen-Shannon Divergence: $JSD(p,q) = \frac{1}{2}KL(p,\frac{p+q}{2}) + \frac{1}{2}KL(q,\frac{p+q}{2})$

$$egin{aligned} &\min \max_{G} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log (1 - D(G(z)))
ight]) \ &= \min_{G} (E_{x \sim p_{data}} \left[\log rac{2 imes p_{data}(x)}{p_{data}(x) + p_G(x)}
ight] + E_{x \sim p_G} \left[\log rac{2 imes p_G(x)}{p_{data}(x) + p_G(x)}
ight] - \log 4) \ &= \min_{G} (KL(p_{data}, rac{p_{data} + p_G}{2}) + KL(p_G, rac{p_{data} + p_G}{2}) - \log 4) \ &= \min_{G} (2 imes JSD(p_{data}, p_G) - \log 4) \end{aligned}$$

Jensen-Shannon Divergence: $JSD(p,q) = \frac{1}{2}KL(p,\frac{p+q}{2}) + \frac{1}{2}KL(q,\frac{p+q}{2})$

$$egin{aligned} &\min \max_{G} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log (1 - D(G(z)))
ight]) \ &= \min_{G} (E_{x \sim p_{data}} \left[\log rac{2 imes p_{data}(x)}{p_{data}(x) + p_{G}(x)}
ight] + E_{x \sim p_{G}} \left[\log rac{2 imes p_{G}(x)}{p_{data}(x) + p_{G}(x)}
ight] - \log 4) \ &= \min_{G} (KL(p_{data}, rac{p_{data} + p_{G}}{2}) + KL(p_{G}, rac{p_{data} + p_{G}}{2}) - \log 4) \ &= \min_{G} (2 imes JSD(p_{data}, p_{G}) - \log 4) \end{aligned}$$

JSD is always nonnegative and zero when the two distributions are equal

=> the global minimum is
$$p_{data}\,=p_G$$

Jensen-Shannon Divergence: $JSD(p,q) = \frac{1}{2}KL(p,\frac{p+q}{2}) + \frac{1}{2}KL(q,\frac{p+q}{2})$

$$\min_{G} \max_{D} (E_{x \sim p_{data}} \left[\log D(x)
ight] + E_{z \sim p(z)} \left[\log(1 - D(G(z)))
ight])$$

$$= \min_G (2*JSD(p_{data},p_G) - \log 4)$$

Summary: The global minimum of the minimax game happens when:

1.
$$D_G^*(x) = rac{p_{data}(x)}{p_{data}(x) + p_G(x)}$$
 (Optimal discriminator for any G)
2. $p_G(x) = p_{data}(x)$ (Optimal generator for optimal D)

(Optimal generator for optimal D)

GAN architecture summary

Jointly train two networks:

Discriminator classifies data as real or fake

Generator generates data that fools the discriminator



Generative Adversarial Networks: results



Generative Adversarial Networks: DC-GAN



Generative Adversarial Networks: Interpolation



Generative Adversarial Networks: Vector Math



Generative Adversarial Networks: Vector Math



Alec Radford, Luke Metz, Soumith Chintala. Unsupervised Representation Learning with Deep Convolutional Generative Adversarial Networks. ICLR 2016

Conditional GANs



[b] Mehdi Mirza, Simon Osindero. Conditional Generative Adversarial Nets. 2014

[c] Augustus Odena, Christopher Olah, Jonathon Shlens. Conditional Image Synthesis With Auxiliary Classifier GANs. ICML 2016

Conditional GANs



monarch butterfly



goldfinch



daisy



redshank



grey whale

Augustus Odena, Christopher Olah, Jonathon Shlens. Conditional Image Synthesis With Auxiliary Classifier GANs. ICML 2016

Conditional GANs: BigGAN



Figure 6: Samples generated by our BigGAN model at 512×512 resolution.



Conditional GANs: Conditional Batch Normalization





The input activation x is normalized across spatial dimensions and scaled and shifted using style-dependent parameter vectors γ_s , β_s where *s* indexes the style label.

Vincent Dumoulin, Jonathon Shlens, Manjunath Kudlur, A Learned Representation For Artistic Style. ICLR 2017 https://arxiv.org/abs/1610.07629

Image Super-Resolution



Christian Ledig, Lucas Theis, Ferenc Huszar, Jose Caballero, Andrew Cunningham, Alejandro Acosta, Andrew Aitken, Alykhan Tejani, Johannes Totz, Zehan Wang, Wenzhe Shi. Photo-Realistic Single Image Super-Resolution Using a Generative Adversarial Network. CVPR 2017 <u>https://arxiv.org/abs/1609.04802</u>

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Image-to-Image Translation: Pix2Pix



Phillip Isola, Jun-Yan Zhu, Tinghui Zhou, Alexei A. Efros. Image-to-Image Translation with Conditional Adversarial Networks. CVPR 2017 https://arxiv.org/abs/1611.07004

Image-to-Image Translation: Pix2Pix

Objective:

 $G^* = \arg\min_{G} \max_{D} \mathcal{L}_{cGAN}(G, D) + \lambda \mathcal{L}_{L1}(G)$

where

$$\mathcal{L}_{cGAN}(G,D) = \mathbb{E}_{x,y}[\log D(x,y)] + \\ \mathbb{E}_{x,z}[\log(1-D(x,G(x,z)))]$$

 $\mathcal{L}_{L1}(G) = \mathbb{E}_{x,y,z}[\|y - G(x,z)\|_1]$



Phillip Isola, Jun-Yan Zhu, Tinghui Zhou, Alexei A. Efros. Image-to-Image Translation with Conditional Adversarial Networks. CVPR 2017 https://arxiv.org/abs/1611.07004



Jun-Yan Zhu, Taesung Park, Phillip Isola, Alexei A. Efros. Unpaired Image-to-Image Translation using Cycle-Consistent Adversarial Networks. ICCV 2017 https://arxiv.org/abs/1703.10593





Jun-Yan Zhu, Taesung Park, Phillip Isola, Alexei A. Efros. Unpaired Image-to-Image Translation using Cycle-Consistent Adversarial Networks. ICCV 2017 https://arxiv.org/abs/1703.10593

Objective:

$$\begin{split} \mathcal{L}(G, F, D_X, D_Y) = & \mathcal{L}_{\text{GAN}}(G, D_Y, X, Y) \\ &+ \mathcal{L}_{\text{GAN}}(F, D_X, Y, X) \\ &+ \lambda \mathcal{L}_{\text{cyc}}(G, F), \end{split}$$

where

$$\mathcal{L}_{\text{GAN}}(G, D_Y, X, Y) = \mathbb{E}_{y \sim p_{\text{data}}(y)} [\log D_Y(y)] + \mathbb{E}_{x \sim p_{\text{data}}(x)} [\log(1 - D_Y(G(x)))],$$

$$\mathcal{L}_{\text{cyc}}(G, F) = \mathbb{E}_{x \sim p_{\text{data}}(x)} [\|F(G(x)) - x\|_1] \\ + \mathbb{E}_{y \sim p_{\text{data}}(y)} [\|G(F(y)) - y\|_1].$$



Jun-Yan Zhu, Taesung Park, Phillip Isola, Alexei A. Efros. Unpaired Image-to-Image Translation using Cycle-Consistent Adversarial Networks. ICCV 2017 https://arxiv.org/abs/1703.10593



Jun-Yan Zhu, Taesung Park, Phillip Isola, Alexei A. Efros. Unpaired Image-to-Image Translation using Cycle-Consistent Adversarial Networks. ICCV 2017 https://arxiv.org/abs/1703.10593

Label Map to Image



Taesung Park, Ming-Yu Liu, Ting-Chun Wang, Jun-Yan Zhu. Semantic Image Synthesis with Spatially-Adaptive Normalization. CVPR 2019 https://arxiv.org/abs/1903.07291

Label Map to Image





Taesung Park, Ming-Yu Liu, Ting-Chun Wang, Jun-Yan Zhu. Semantic Image Synthesis with Spatially-Adaptive Normalization. CVPR 2019 https://arxiv.org/abs/1903.07291
StyleGAN





Tero Karras, Samuli Laine, Timo Aila. A Style-Based Generator Architecture for Generative Adversarial Networks. CVPR 2019

Video Generation



Aliaksandr Siarohin, Stéphane Lathuilière, Sergey Tulyakov, Elisa Ricci, Nicu Sebe. First Order Motion Model for Image Animation. NeurIPS 2019

Video Generation



Video Generation. Everybody Dance Now



Caroline Chan, Shiry Ginosar, Tinghui Zhou, Alexei A. Efros. Everybody Dance Now. ICCV 2019

Video Generation. Everybody Dance Now



Caroline Chan, Shiry Ginosar, Tinghui Zhou, Alexei A. Efros. Everybody Dance Now. ICCV 2019

Problems of GANs

- Mode collapse:
 - G collapses providing limited sample variety
- Non-convergence:
 - model parameters oscillate, destabilize and never converge
- Diminished gradient:
 - D is so successful that the G gradient vanishes and learns nothing

Mode collapse

Real-life data is multimodal (10 in MNIST)

Mode collapse: when few modes generated



Partial mode collapse

The generator produces realistic and diverse samples, but much less diverse than the real-world data distribution.



Solutions to mode collapse

- Wasserstein loss [1]
 - Trains the discriminator to optimality without worrying about vanishing gradients.
 - If the discriminator doesn't get stuck in local minima, it learns to reject the outputs that the generator stabilizes on.
- Unrolling [2]
 - Uses a generator loss function that incorporates not only the current discriminator's classifications, but also the outputs of future discriminator versions
 - The generator can't over-optimize for a single discriminator.

[1] Martin Arjovsky, Soumith Chintala, Léon Bottou. Wasserstein GAN. 2017 https://arxiv.org/abs/1701.07875

[2] Luke Metz, Ben Poole, David Pfau, Jascha Sohl-Dickstein. Unrolled Generative Adversarial Networks. https://arxiv.org/abs/1611.02163

• GAN can optimize the discriminator easier than the generator.

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Original GAN generator's gradient:
$$-\nabla_{\theta_g} \log \left(1 - D\left(G\left(\boldsymbol{z}^{(i)}\right)\right)\right) \rightarrow \boldsymbol{\theta}$$

Alternative: $\nabla_{\theta_g} \log D\left(G\left(\boldsymbol{z}^{(i)}\right)\right)$

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Alternative:
$$\nabla_{\theta_g} \log D\left(G\left(\boldsymbol{z}^{(i)}\right)\right)$$

Problem: large variance of gradients that make the model unstable



Martin Arjovsky, Soumith Chintala, Léon Bottou. Wasserstein GAN. 2017 https://arxiv.org/abs/1701.07875

$$W(\mathbb{P}_r, \mathbb{P}_{\theta}) = \sup_{\|f\|_L \le 1} \mathbb{E}_{x \sim \mathbb{P}_r}[f(x)] - \mathbb{E}_{x \sim \mathbb{P}_{\theta}}[f(x)]$$



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where

- *sup* is the least upper bound
- *f* is a 1-Lipschitz function following constraint:

$$|f(x_1)-f(x_2)|\leq |x_1-x_2|.$$



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We can build a deep network to calculate the Wasserstein distance.



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We can build a deep network to calculate the Wasserstein distance.

This network is very similar to the discriminator, just without the sigmoid function and outputs a scalar score rather than a probability.



$$W(\mathbb{P}_r, \mathbb{P}_{\theta}) = \sup_{\|f\|_L \le 1} \mathbb{E}_{x \sim \mathbb{P}_r}[f(x)] - \mathbb{E}_{x \sim \mathbb{P}_{\theta}}[f(x)]$$

Algorithm 1 WGAN, our proposed algorithm. All experiments in the paper used the default values $\alpha = 0.00005$, c = 0.01, m = 64, $n_{\text{critic}} = 5$.

Require: : α, the learning rate. c, the clipping parameter. m, the batch size. n_{critic}, the number of iterations of the critic per generator iteration.
Require: : w₀, initial critic parameters. θ₀, initial generator's parameters.
1: while θ has not converged do
2: for t = 0, ..., n_{critic} do

3: Sample
$$\{x^{(i)}\}_{i=1}^{m} \sim \mathbb{P}_{r}$$
 a batch from the real data.
4: Sample $\{z^{(i)}\}_{i=1}^{m} \sim p(z)$ a batch of prior samples.
5: $g_{w} \leftarrow \nabla_{w} \left[\frac{1}{m} \sum_{i=1}^{m} f_{w}(x^{(i)}) - \frac{1}{m} \sum_{i=1}^{m} f_{w}(g_{\theta}(z^{(i)}))\right]$
6: $w \leftarrow w + \alpha \cdot \text{RMSProp}(w, g_{w})$
7: $w \leftarrow \text{clip}(w, -c, c)$
8: **end for**
9: Sample $\{z^{(i)}\}_{i=1}^{m} \sim p(z)$ a batch of prior samples.
10: $g_{\theta} \leftarrow -\nabla_{\theta} \frac{1}{m} \sum_{i=1}^{m} f_{w}(g_{\theta}(z^{(i)}))$
11: $\theta \leftarrow \theta - \alpha \cdot \text{RMSProp}(\theta, g_{\theta})$
12: **end while**





Martin Arjovsky, Soumith Chintala, Léon Bottou. Wasserstein GAN. 2017 https://arxiv.org/abs/1701.07875

• Wasserstein criterion allows us to train **D** until optimality. When the criterion reaches the optimal value, it simply provides a loss to the generator that we can train as any other neural network.

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- Wasserstein criterion allows us to train **D** until optimality. When the criterion reaches the optimal value, it simply provides a loss to the generator that we can train as any other neural network.
- We no longer need to balance **G** and **D** capacity properly.
- Wasserstein loss leads to a higher quality of the gradients to train **G**.
- WGANs are **more robust** than common GANs to the architectural choices for the generator and hyperparameter tuning

Evaluation

The objective function for the generator and the discriminator usually measures how well they are doing relative to the opponent.

It is not a good metric in measuring the image quality or its diversity.

Evaluation

- Inception Score (IS) [1]
- Frechet Inception Distance (FID) [2]
- Human-based ratings and preference judgments

[1] Tim Salimans, Ian Goodfellow, Wojciech Zaremba, Vicki Cheung, Alec Radford, Xi Chen. Improved Techniques for Training GANs. NeurIPS 2016 https://arxiv.org/abs/1606.03498

[2] Martin Heusel, Hubert Ramsauer, Thomas Unterthiner, Bernhard Nessler, Sepp Hochreiter. GANs Trained by a Two Time-Scale Update Rule Converge to a Local Nash Equilibrium. NeurIPS 2017 <u>https://arxiv.org/abs/1706.08500</u>

IS uses two criteria in measuring the performance of GAN:

- The **quality** of the generated images
- their **diversity**

- Quality: use an Inception network to predict conditional probability p(y|x) where y is the label and x is the generated $p(y) = \int_{z}^{z} p(y|x = G(z)) dz$
- **Diversity**: calculate marginal probability:

- **Quality:** use an Inception network to predict conditional probability p(y|x) where y is the label and x is the generated data **Diversity**: calculate marginal probability: $p(y) = \int_z p(y|x = G(z))dz$

We want

- the conditional probability p(y|x) to be highly predictable (low entropy) i.e. given an image, we should know the object type easily
- the data distribution **p(y)** should be uniform (high entropy)

P(y)

```
Inception Score (IS)
```

Compute their KL-divergence to combine these two criteria: $IS(G) = \exp(E_{x \sim p_g} KL(p(y|x)||p(y)))$

Limitations:

- IS is limited by what the Inception classifier can detect, which is linked to the training data (ILSVRC)
- IS can misrepresent the performance if it only generates one image per class.
 p(y) will still be uniform even though the diversity is low

• Use the Inception network to extract features from an intermediate layer

Martin Heusel, Hubert Ramsauer, Thomas Unterthiner, Bernhard Nessler, Sepp Hochreiter. GANs Trained by a Two Time-Scale Update Rule Converge to a Local Nash Equilibrium. NeurIPS 2017 <u>https://arxiv.org/abs/1706.08500</u>

- Use the Inception network to extract features from an intermediate layer
- Model data distribution for these features using a multivariate Gaussian distribution with mean μ and covariance Σ

- Use the Inception network to extract features from an intermediate layer
- Model data distribution for these features using a multivariate Gaussian distribution with mean μ and covariance Σ
- The FID between the real images x and generated images g: $FID(x,g) = ||\mu_x - \mu_g||_2^2 + Tr(\Sigma_x + \Sigma_g - 2(\Sigma_x \Sigma_g)^{\frac{1}{2}})$

where Tr sums up all the diagonal elements

- Lower FID values mean better image quality and diversity
- FID is sensitive to mode collapse, the distance increases when modes are missed
- FID is more robust to noise than IS. If the model only generates one image per class, the distance will be high

