Deep Generative Learning

Deep Learning for Computer Vision

Valeriya Strizhkova

21.11.23

About myself

Valeriya Strizhkova

PhD candidate @ Inria & 3iA & UCA

Research topics:

- video understanding
- human behavior understanding
- multi-sensor fusion

valeriya.strizhkova@inria.fr



Outline

Introduction

- What is a generative model?
- Types of deep generative models
- Evaluation: IS, FID

Image Generation

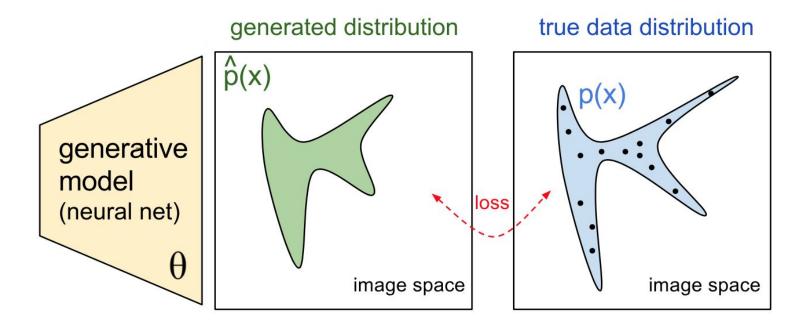
- Autoregressive models: PixelCNN, DALL-E 1
- Variational Autoencoders: dVAE
- Generative Adversarial Networks: StyleGAN,
- Diffusion Models: Latent diffusion, DALL-E 2

Theory

- Variational lower bound maximization
- GAN training objective and optimality

Introduction

What is a Generative Model?



The Landscape of Deep Generative Learning

Autoregressive Models Normalizing Flows

Variational Autoencoders

Generative Adversarial Networks

Energy-based Models Denoising
Diffusion Models

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

Types of generative models

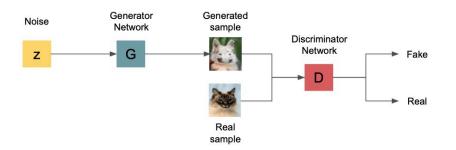
Variational Autoencoder:

mean vector sampled latent vector Encoder Network (conv) Standard deviation vector

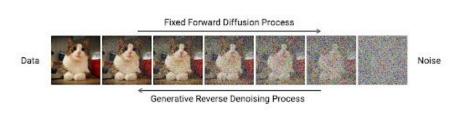
Autoregressive model:



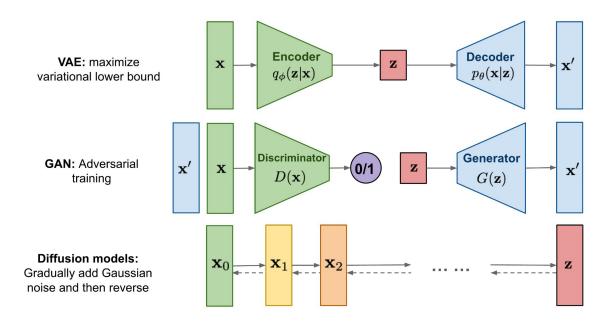
Generative Adversarial Network:



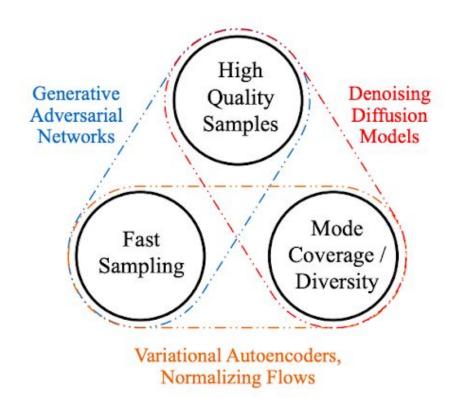
Diffusion model:



Types of generative models



Generative learning trilemma

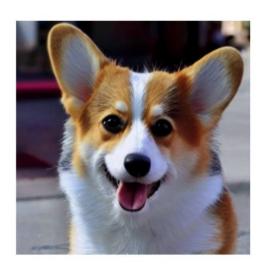


Which of these images looks better?





Which of these images looks more realistic?





Which of these images appears to be more similar to the text prompt?





Prompt: The saying "BE EXCELLENT TO EACH OTHER" written on a red brick wall with a graffiti image of a green alien wearing a tuxedo. A yellow fire hydrant is on a sidewalk in the foreground.

- Human-based ratings and preference judgments
- Inception Score (IS) [1]
- Frechet Inception Distance (FID) [2]
- Structured Similarity Index Metric (SSIM) [3]
- CLIP score [4]

[4] Jack Hessel, Ari Holtzman, Maxwell Forbes, Ronan Le Bras, Yejin Choi.CLIPScore: A Reference-free Evaluation Metric for Image Captioning. EMNLP 2021https://arxiv.org/abs/2104.08718

^[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

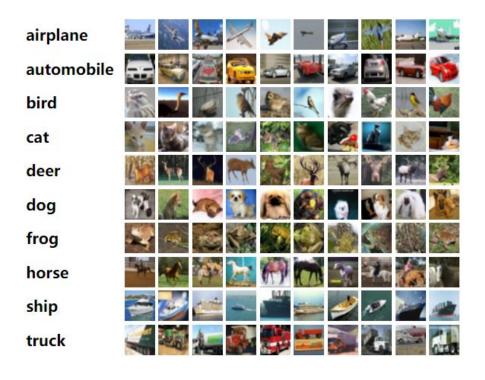
^[2] Martin Heusel, Hubert Ramsauer, Thomas Unterthiner, Bernhard Nessler, Sepp Hochreiter. GANs Trained by a Two Time-Scale Update Rule Converge to a Local Nash Equilibrium. NeurIPS 2017 https://arxiv.org/abs/1706.08500

^[3] Zhou Wang; A.C. Bovik; H.R. Sheikh; E.P. Simoncelli. Image quality assessment: from error visibility to structural similarity. IEEE Transactions on Image Processing 2004 https://www.cns.nyu.edu/pub/eero/wang03-reprint.pdf

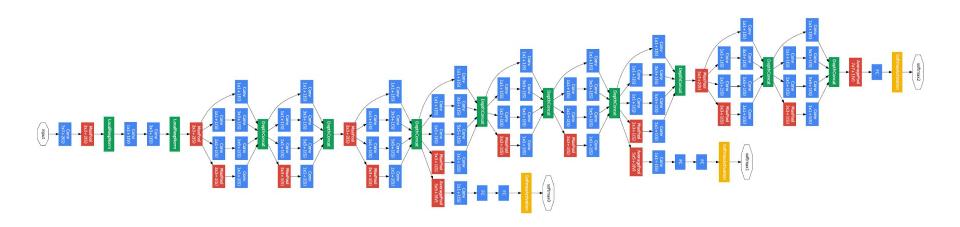
IS measures:

- the quality of the generated images
- their diversity

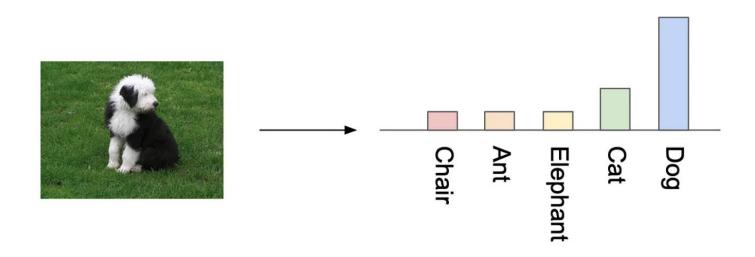
CIFAR10 dataset:



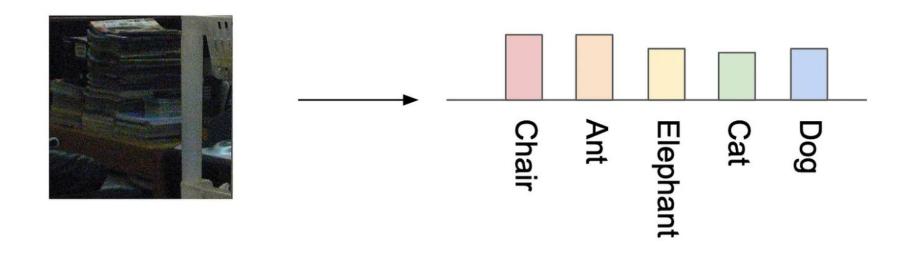
Generate 50000 images by the **Inception** image classifier pre-trained on CIFAR10



Generated images are fed into the Inception image classifier network pre-trained on the CIFAR10 dataset predict conditional probability p(y|x) — where y is the label and x is the generated data

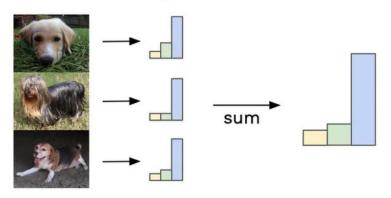


If the probability scores are widely distributed then the generated image is of low quality:

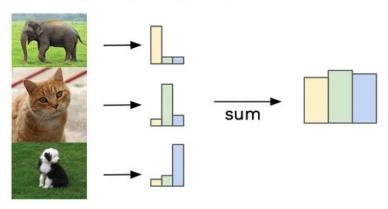


Calculate marginal probability $p(y) = \int_z p(y|x=G(z))dz$

Similar labels sum to give focussed distribution



Different labels sum to give uniform distribution



- Quality: conditional probability p(y|x)
- Diversity: marginal probability p(y)

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 marginal probability p(y) to be uniform (high entropy)

Compute their KL-divergence to combine these two criteria:

$$IS(G) = \exp(E_{x \sim p_q} KL(p(y|x)||p(y)))$$

Inception Score (IS) limitations

IS is limited by what the Inception classifier can detect, which is linked to the training data

FID compares the distribution of generated images with the distribution of real images

Uses the Inception network to extract features from an intermediate layer

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

- Uses the Inception network to extract features from an intermediate layer
- Models 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 is:

$$|FID(x,g) = \left|\left|\mu_x - \mu_g
ight|
ight|^2_2 + Tr(\Sigma_x + \Sigma_g - 2(\Sigma_x\Sigma_g)^{rac{1}{2}})
ight|$$

where Tr sums up all the diagonal elements

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

Structured Similarity Index Metric (SSIM)

The Structural Similarity Index (SSIM) metric mimicks the human visual perception system which is highly capable of identifying **structural information** from a scene.

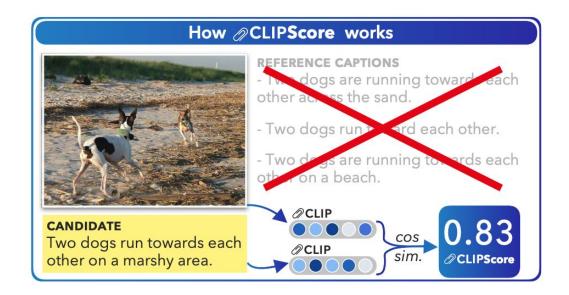
SSIM extracts 3 features from an image on a **pixel-wise** basis:

- Luminance
- Contrast
- Structure

The comparison between the two images is performed on the basis of these 3 features.

CLIP score

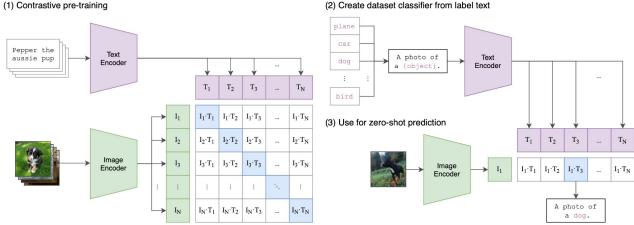
- CLIP score measures the compatibility of image-caption pairs.
- Higher CLIP scores imply higher compatibility



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

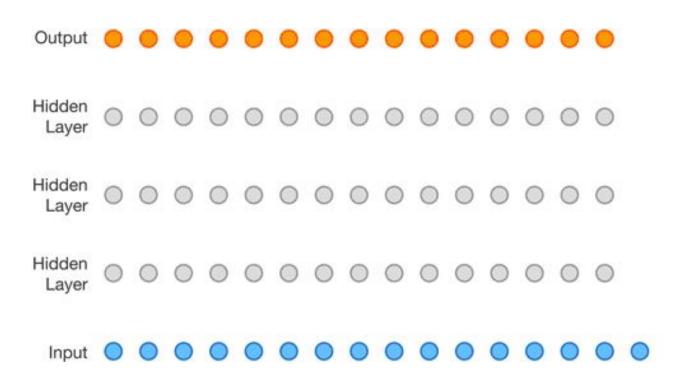


CLIP Architecture

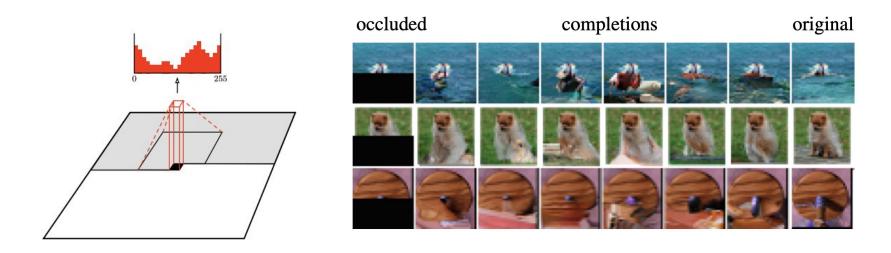
- text and image have separate transformer encoders
- visual encoder is ViT (vision transformer)
- text encoder is GPT-2 transformer
- the fixed-length text embedding is extracted from [EOS] token position
- text token embeddings and image patch embeddings also available
- trained on 256 GPUs for 2 weeks

Image Generation

Autoregressive Models



Pixel Recurrent Neural Networks



DALL-E 1

- is introduced by OpenAl
- generates 256×256 images from text via dVAE
- autoregressively generates image tokens from textual tokens on a discrete latent space

DALL-E 1 Training

- 1. 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

- encode input 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

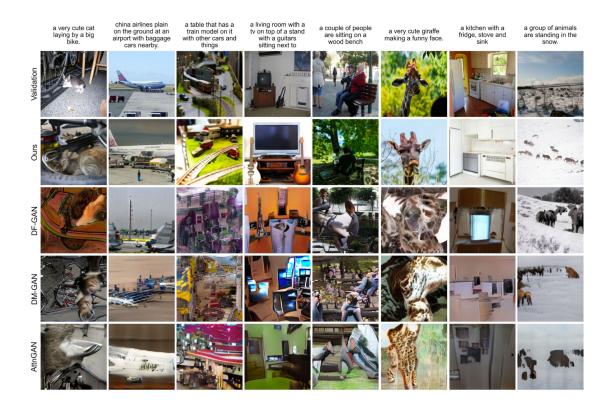


accordion.

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

DALL-E 1 Results



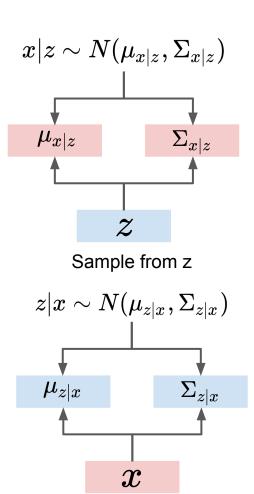
Aditya Ramesh, Mikhail Pavlov, Gabriel Goh, Scott Gray, Chelsea Voss, Alec Radford, Mark Chen, Ilya Sutskever. Zero-Shot Text-to-Image Generation. PMLR 2021 https://arxiv.org/abs/2102.12092

Variational Autoencoders (VAE)

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)
- A point from the latent space is sampled from that distribution
- 4. The sampled point is **decoded**
- 5. The reconstruction error is computed



Variational Autoencoders (VAE)

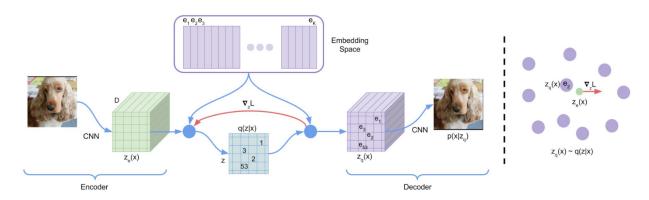
Learned data manifold for generative models with two-dimensional latent space:



```
000000000000000
```

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

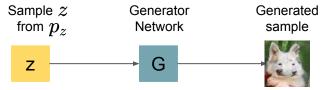
• Setup: Assume we have data x_i drawn from distribution $p_{data}(x)$. Want to sample from p_{data} .

- Setup: Assume we have data x_i drawn from distribution $p_{data}(x)$. Want to sample from p_{data} .
- Idea: Introduce a latent variable z with simple prior p(z).

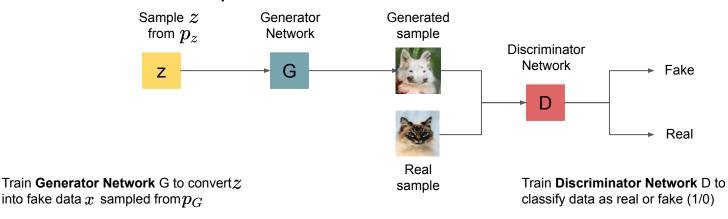
- Setup: Assume we have data x_i drawn from distribution $p_{data}(x)$. Want to sample from p_{data} .
- Idea: Introduce a latent variable z with simple prior p(z).
- ullet Sample $z\sim p(z)$ and pass to a Generator Network $\,x=G(z)\,$

- Setup: Assume we have data x_i drawn from distribution $p_{data}(x)$. Want to sample from p_{data} .
- Idea: Introduce a latent variable z with simple prior p(z).
- ullet Sample $z\sim p(z)$ and pass to a Generator Network $\,x=G(z)\,$
- Then x is a sample from the Generator distribution p_G . Want $p_G = p_{data}$

- Setup: Assume we have data x_i drawn from distribution $p_{data}(x)$. Want to sample from p_{data} .
- Idea: Introduce a latent variable z with simple prior p(z).
- ullet Sample $z\sim p(z)$ and pass to a Generator Network $\,x=G(z)\,$
- ullet Then x is a sample from the Generator distribution $\,p_G$. Want $\,p_G=p_{data}$

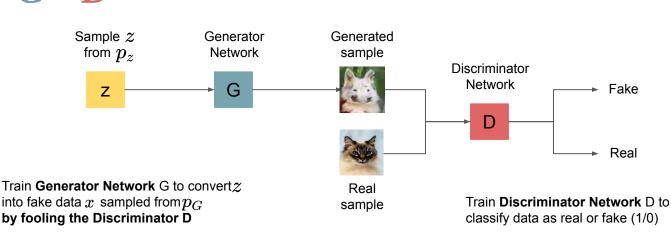


- Setup: Assume we have data x_i drawn from distribution $p_{data}(x)$. Want to sample from p_{data} .
- Idea: Introduce a latent variable z with simple prior p(z).
- ullet Sample $z\sim p(z)$ and pass to a Generator Network $\,x=G(z)\,$
- ullet Then x is a sample from the Generator distribution p_G . Want $p_G=p_{data}$



Jointly train generator G and discriminator D with a minimax game

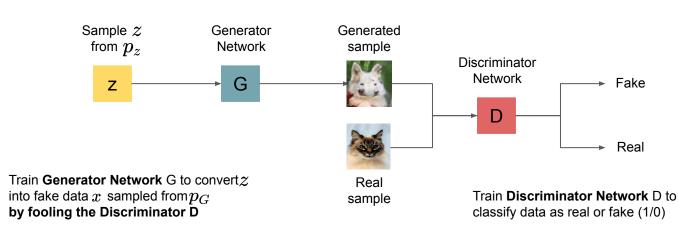
$$\min \max_{\mathbf{G}} (E_{x \sim p_{data}}[\log \mathbf{D}(x)] + E_{\mathbf{Z} \sim p(\mathbf{Z})}[\log(1 - \mathbf{D}(\mathbf{G}(\mathbf{Z})))])$$



Jointly train generator G and discriminator D with a minimax game

Discriminator wants D(x)=1 for real data

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

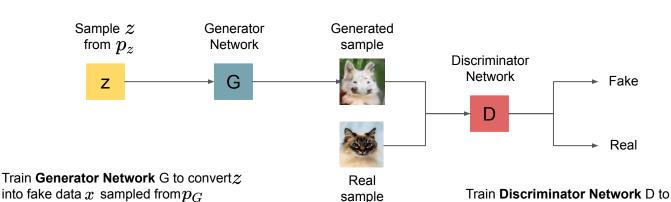


Jointly train generator G and discriminator D with a minimax game

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

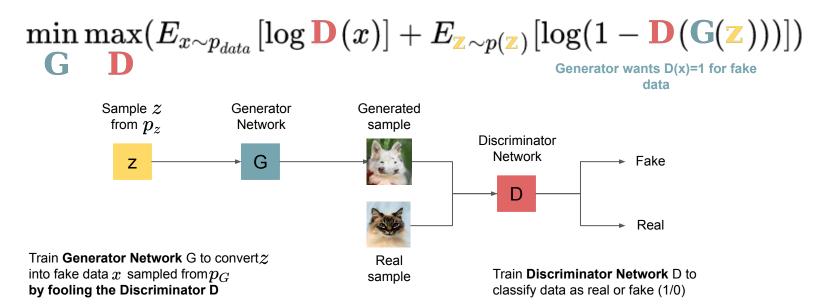
Discriminator wants D(x)=0 for

classify data as real or fake (1/0)



by fooling the Discriminator D

Jointly train generator G and discriminator D with a minimax game



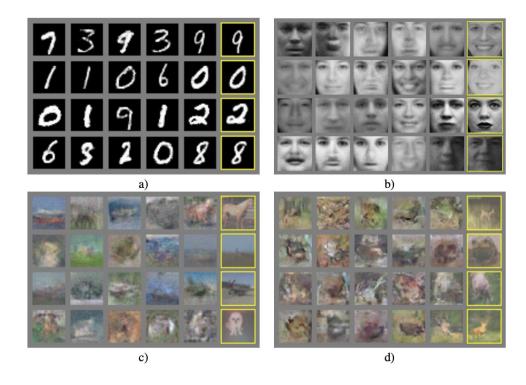
Jointly train generator G and discriminator D with a minimax game

$$egin{aligned} \min \max_{\mathbf{G}} \left(E_{x \sim p_{data}} \left[\log \mathbf{D}(x) \right] + E_{\mathbf{Z} \sim p(\mathbf{Z})} \left[\log (1 - \mathbf{D}(\mathbf{G}(\mathbf{Z}))) \right]
ight) \ &= \min \max_{\mathbf{G}} \mathbf{V}(\mathbf{G}, \mathbf{D}) \end{aligned}$$

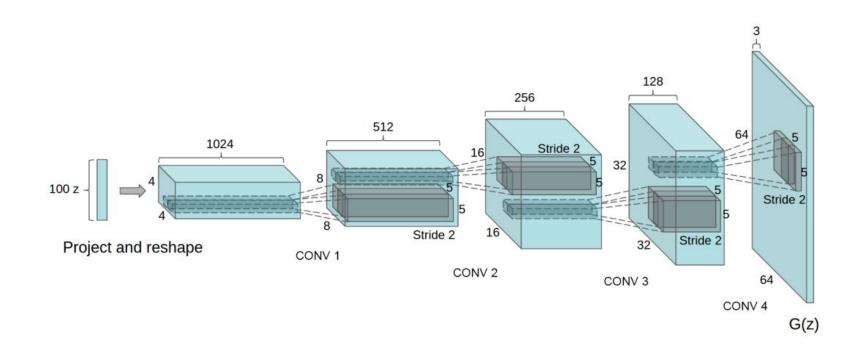
Train G and D using alternating gradient updates:

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: results



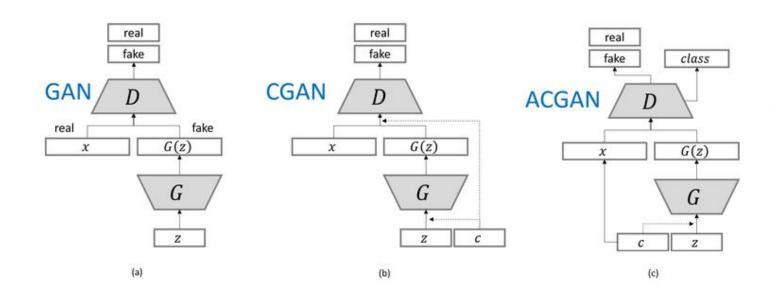
Generative Adversarial Networks: DC-GAN



Generative Adversarial Networks: Interpolation



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

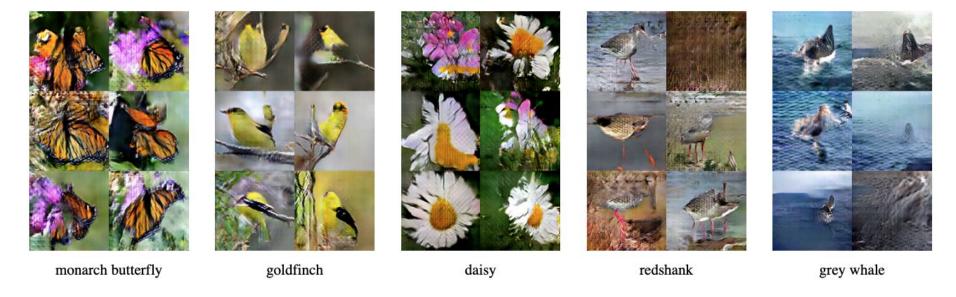


Image-to-Image Translation: Pix2Pix

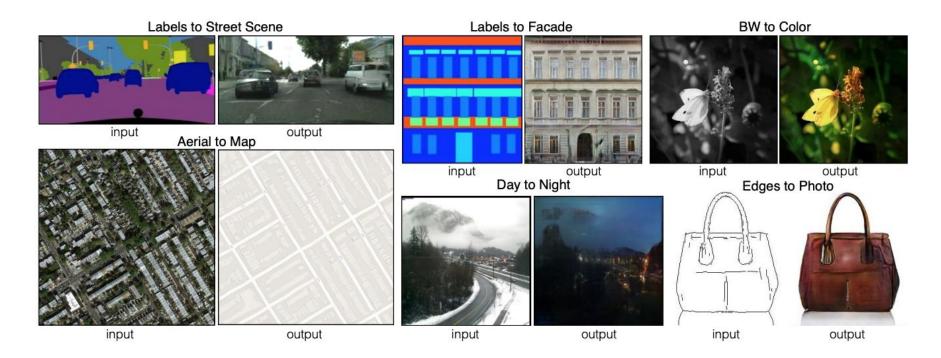


Image-to-Image Translation: Pix2Pix

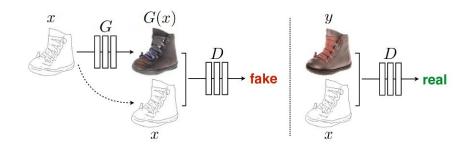
Objective:

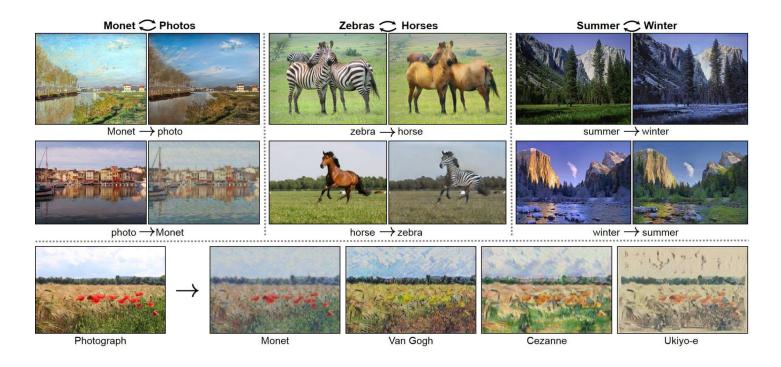
$$G^* = \arg\min_{G} \max_{D} \mathcal{L}_{cGAN}(G, D) + \lambda \mathcal{L}_{L1}(G)$$

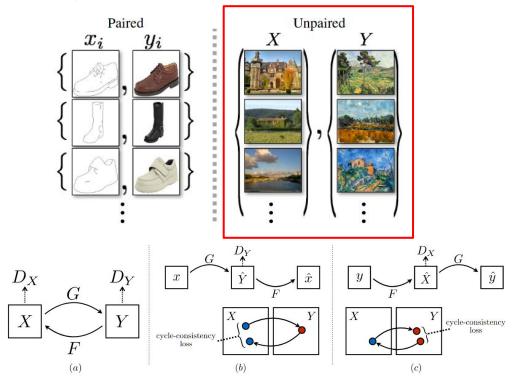
where

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

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







Objective:

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

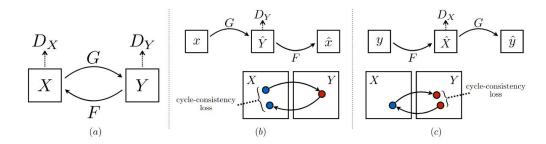
where

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

$$+ \mathbb{E}_{x \sim p_{\text{data}}(x)}[\log(1 - D_Y(G(x))],$$

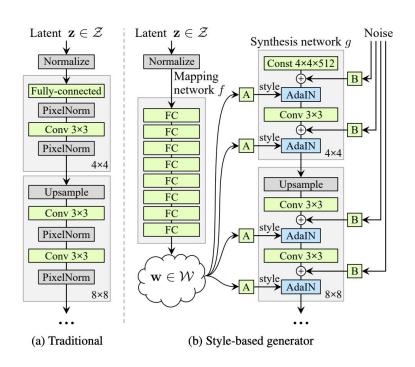
$$\mathcal{L}_{\text{cyc}}(G, F) = \mathbb{E}_{x \sim p_{\text{data}}(x)}[\|F(G(x)) - x\|_1]$$

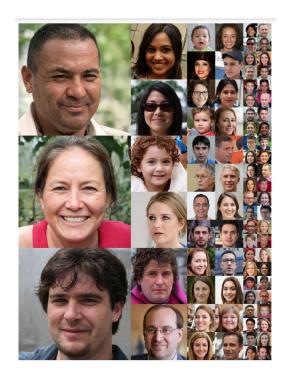
$$+ \mathbb{E}_{y \sim p_{\text{data}}(y)}[\|G(F(y)) - y\|_1].$$





StyleGAN





StyleGAN



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

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, 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:

Data

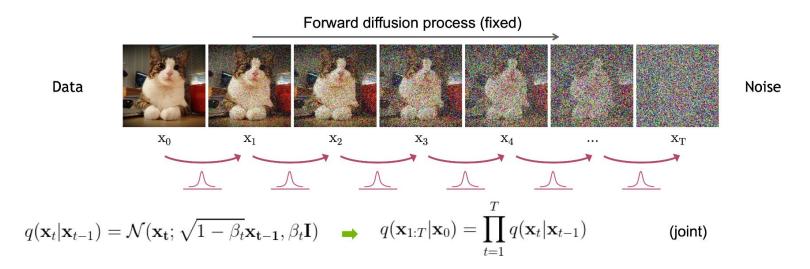
- Forward diffusion process that gradually adds noise to input
- Reverse denoising process that learns to generate data by denoising

Forward diffusion process (fixed) 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

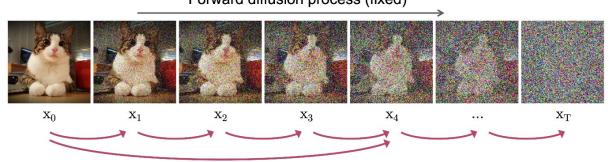
The formal definition of the forward process in T steps:



Diffusion Kernel



Data



Noise

Define
$$\bar{\alpha}_t = \prod_{s=1}^t (1 - \beta_s)$$
 \Rightarrow $q(\mathbf{x}_t | \mathbf{x}_0) = \mathcal{N}(\mathbf{x}_t; \sqrt{\bar{\alpha}_t} \mathbf{x}_0, (1 - \bar{\alpha}_t) \mathbf{I}))$ (Diffusion Kernel) For sampling: $\mathbf{x}_t = \sqrt{\bar{\alpha}_t} \ \mathbf{x}_0 + \sqrt{(1 - \bar{\alpha}_t)} \ \epsilon$ where $\epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$

For sampling:
$$\mathbf{x}_t = \sqrt{\bar{\alpha}_t} \; \mathbf{x}_0 + \sqrt{(1 - \bar{\alpha}_t)} \; \epsilon$$
 where $\epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$

 eta_t values schedule (i.e., the noise schedule) is designed such that $ar{lpha}_T o 0$ and $q(\mathbf{x}_T|\mathbf{x}_0) pprox \mathcal{N}(\mathbf{x}_T;\mathbf{0},\mathbf{I}))$

Learning Denoising Model

Variational upper bound

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

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

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 $\tilde{\beta}_t := \frac{1 - \bar{\alpha}_{t-1}}{1 - \bar{\alpha}_t} \beta_t$

Summary

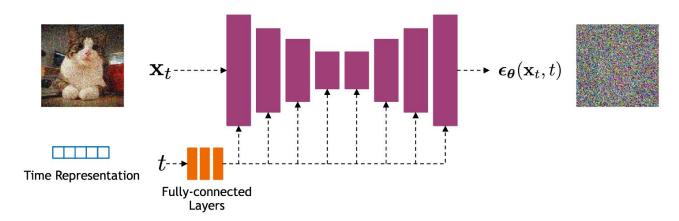
Training and Sample Generation

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: $\epsilon \sim \mathcal{N}(0, \mathbf{I})$ 5: Take gradient descent step on $\nabla_{\theta} \ \epsilon - \epsilon_{\theta} (\sqrt{\bar{\alpha}_t} \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \epsilon, t) \ ^2$ 6: until converged	1: $\mathbf{x}_{T} \sim \mathcal{N}(0, \mathbf{I})$ 2: for $t = T, \dots, 1$ do 3: $\mathbf{z} \sim \mathcal{N}(0, \mathbf{I})$ 4: $\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_{t}}} \left(\mathbf{x}_{t} - \frac{1-\alpha_{t}}{\sqrt{1-\bar{\alpha}_{t}}} \boldsymbol{\epsilon}_{\theta}(\mathbf{x}_{t}, t) \right) + \sigma_{t} \mathbf{z}$ 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 $\epsilon_{\theta}(\mathbf{x}_t,t)$

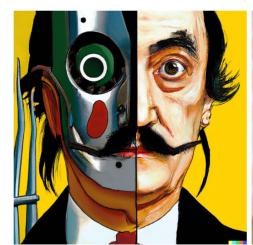


Time representation: sinusoidal positional embeddings or random Fourier features.

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

DALL-E 2

- is introduced by OpenAl
- 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



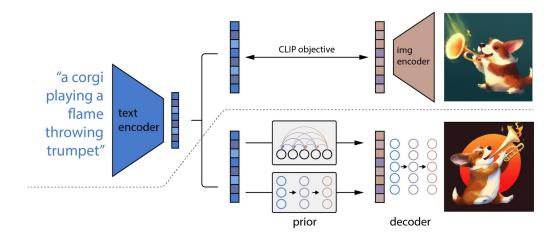
a close up of a handpalm with leaves growing from it

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



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 training

Can vary images while preserving style and semantics in the embeddings



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.

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 Limitations



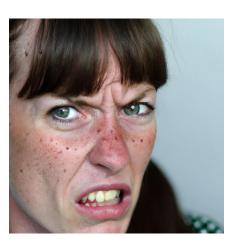
A family dining for Christmas



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

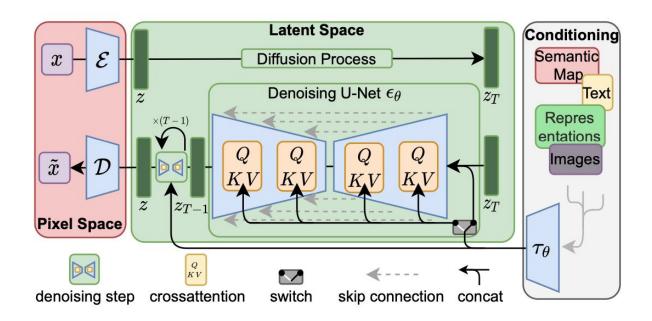




Robin Rombach, Andreas Blattmann, Dominik Lorenz, Patrick Esser, Björn Ommer.

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

Latent Diffusion Models



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



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









Credits to

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

https://openai.com/research/generative-models

https://lilianweng.github.io/posts/2021-07-11-diffusion-models/

https://developer.nvidia.com/blog/improving-diffusion-models-as-an-alternative-to-gans-part-1/