

# CS 4803-DL / 7643-A: Lecture 19

## Danfei Xu

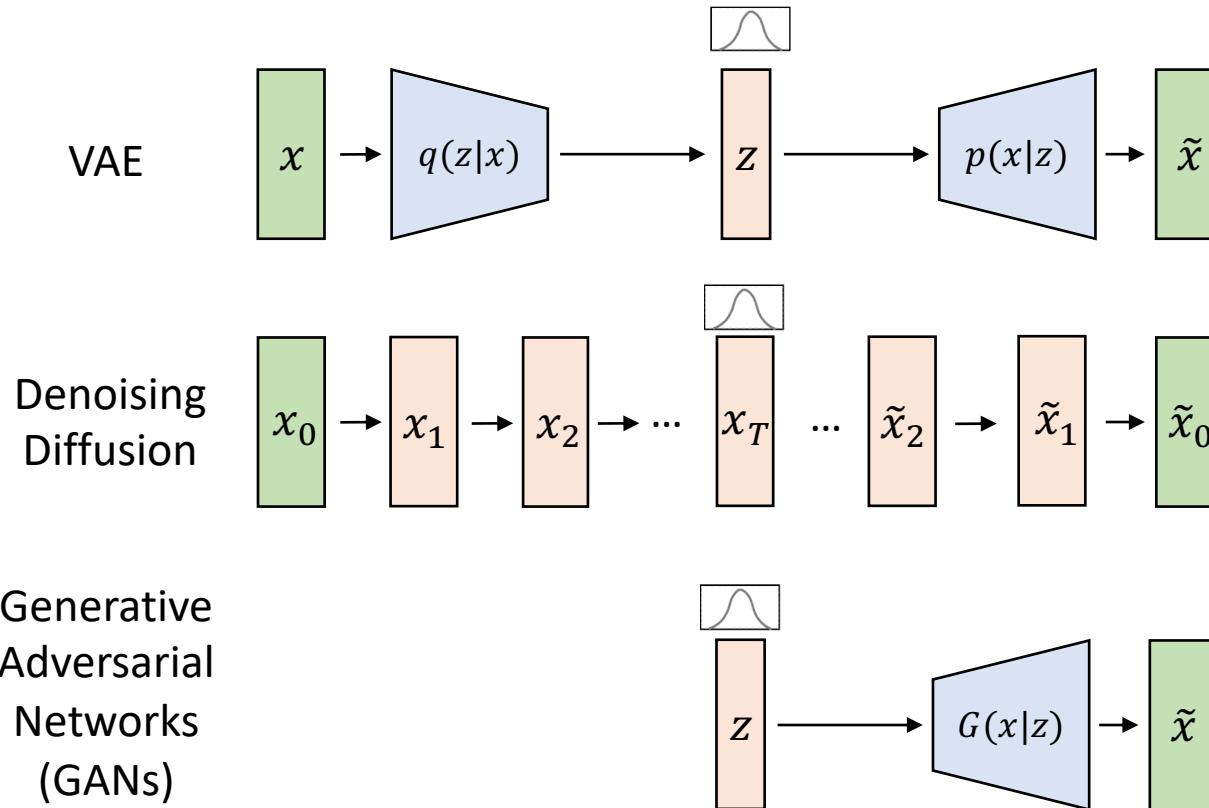
Topics:

- Generative Adversarial Networks
- Self-supervised Learning
  - Pretext task from image transformation
  - Contrastive learning

# Administrative

- HW2 / PS2 grade out. Please submit your regrade request by the end of this week
- HW4 / PS4 out. Due Nov 14<sup>th</sup>. Grade Period ends 16<sup>th</sup>.
- **Start the coding part NOW** --- it takes some time to run GAN / diffusion model training on Colab GPUs.
- Milestone Report & Video due Nov 7<sup>th</sup>. **NO GRACE PERIOD**

# Denoising Diffusion: Image to Noise and Back

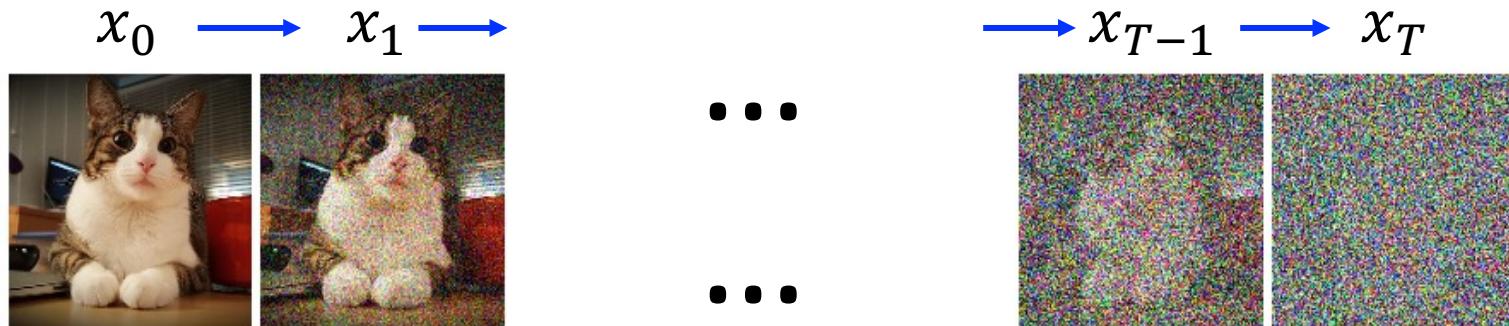


# The Denoising Diffusion Process

image from  
dataset

The “forward diffusion” process:  
add Gaussian noise each step

noise  $\mathcal{N}(0, I)$



$x_0 \leftarrow x_1 \leftarrow \dots \leftarrow x_{T-1} \leftarrow x_T$

The “denoising diffusion” process:  
generate an image from noise by  
*denoising* the gaussian noises

# The Denoising (Decoding) Process

The **learned** denoising process  $x_0 \leftarrow x_1 \leftarrow \dots \leftarrow x_T$

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

**High-level intuition:** derive a *ground truth denoising distribution*  $q(x_{t-1}|x_t, x_0)$  and train a neural net  $p_\theta(x_{t-1}|x_t)$  to match the distribution.

**The learning objective:**  $\text{argmin}_\theta D_{KL}(q(x_{t-1}|x_t, x_0) || p_\theta(x_{t-1}|x_t))$

What does it look like?  $q(x_{t-1}|x_t, x_0) = \mathcal{N}\left(x_{t-1}; \mu_q(t), \Sigma_q(t)\right)$

$$\mu_q(t) = \frac{1}{\sqrt{\alpha_t}} \left( x_t - \frac{\beta_t}{\sqrt{(1 - \bar{\alpha}_t)}} \epsilon \right), \quad \epsilon \sim \mathcal{N}(0, I) \leftarrow \begin{array}{l} \text{Recall: Gaussian} \\ \text{reparameterization trick} \end{array}$$

The “ground truth” noise that brought  $x_{t-1}$  to  $x_t$

# The Denoising (Decoding) Process

The **learned** denoising process  $x_0 \leftarrow x_1 \leftarrow \dots \leftarrow x_T$

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What does it look like?  $q(x_{t-1}|x_t, x_0) = \mathcal{N}(x_{t-1}; \mu_q(t), \Sigma_q(t))$

Assuming identical variance  $\Sigma_q(t)$ , we have:

$$\text{argmin}_\theta D_{KL}(q(x_{t-1}|x_t, x_0) || p_\theta(x_{t-1}|x_t)) = \text{argmin}_\theta w \|\mu_q(t) - \mu_\theta(x_t, t)\|$$

Should be variance-dependent, but constant  
works better in practice

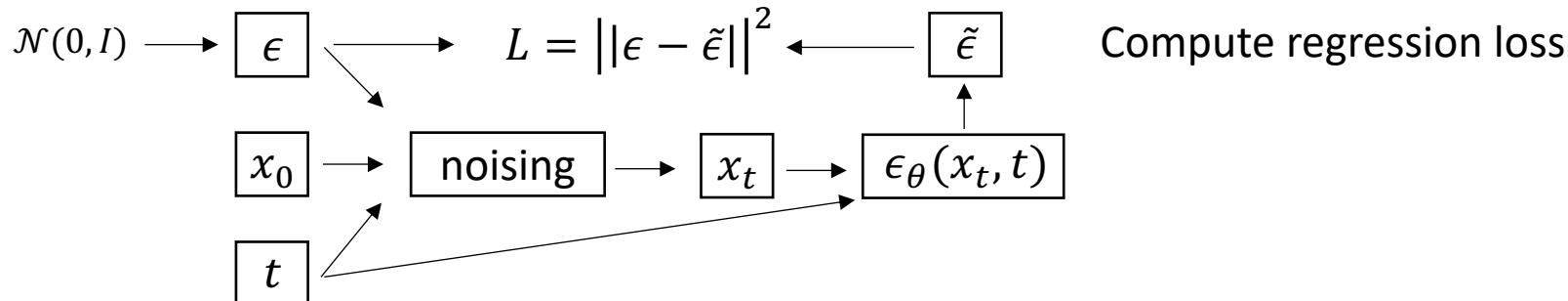
# The Denoising Diffusion Algorithm

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## Algorithm 1 Training

---

- 1: **repeat**
  - 2:  $\mathbf{x}_0 \sim q(\mathbf{x}_0)$
  - 3:  $t \sim \text{Uniform}(\{1, \dots, T\})$
  - 4:  $\epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$
  - 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
- 



# The Denoising Diffusion Algorithm

---

## Algorithm 1 Training

---

```
1: repeat
2:    $\mathbf{x}_0 \sim q(\mathbf{x}_0)$ 
3:    $t \sim \text{Uniform}(\{1, \dots, T\})$ 
4:    $\epsilon \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ 
5:   Take gradient descent step on
      
$$\nabla_{\theta} \left\| \epsilon - \epsilon_{\theta}(\sqrt{\bar{\alpha}_t} \mathbf{x}_0 + \sqrt{1 - \bar{\alpha}_t} \epsilon, t) \right\|^2$$

6: until converged
```

---

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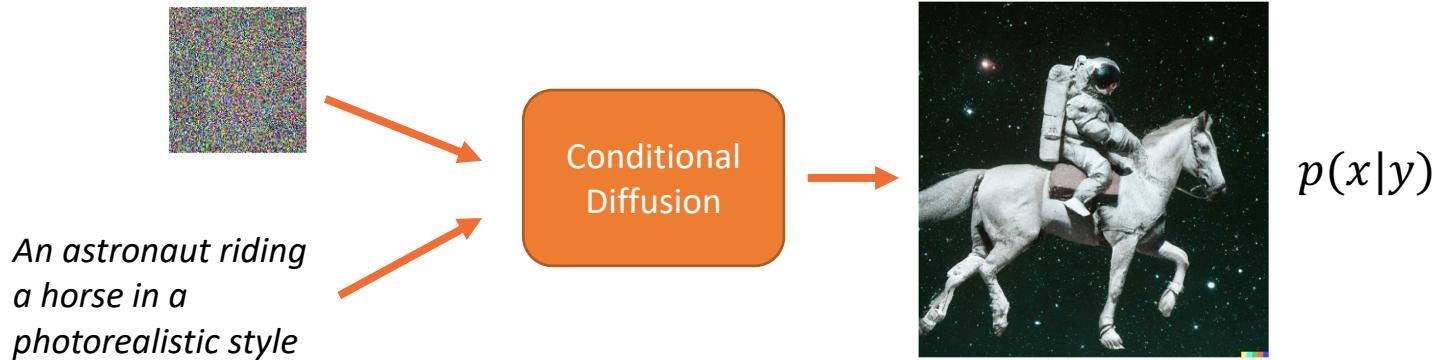
## Algorithm 2 Sampling

---

```
1:  $\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ 
2: for  $t = T, \dots, 1$  do
3:    $\mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$  if  $t > 1$ , else  $\mathbf{z} = \mathbf{0}$ 
4:    $\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left( \mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1 - \bar{\alpha}_t}} \epsilon_{\theta}(\mathbf{x}_t, t) \right) + \sigma_t \mathbf{z}$ 
5: end for
6: return  $\mathbf{x}_0$ 
```

---

# Classifier-free Guided Diffusion

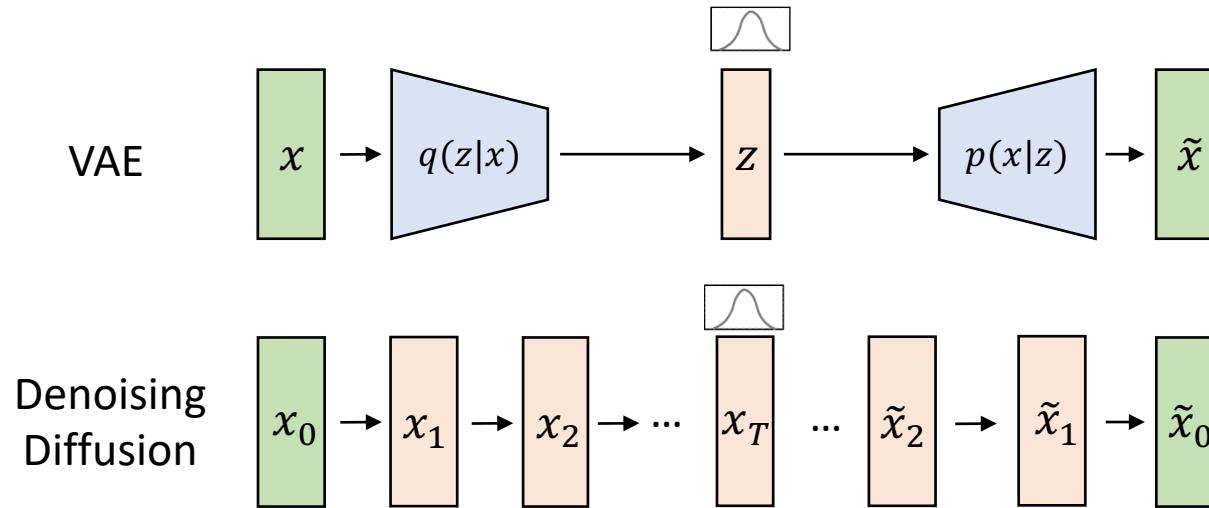


**Classifier-free Guided Diffusion:** estimate the gradient of the classifier model with conditional diffusion models!

$$\nabla_{x_t} \log f_\varphi(y|x_t) = -\frac{1}{\sqrt{1-\bar{\alpha}_t}} (\epsilon_\theta(x_t, t, y) - \bar{\epsilon}_\theta(x_t, t))$$
$$\bar{\epsilon}_\theta(x_t, t, y) = (w+1)\epsilon_\theta(x_t, t, y) - w\epsilon_\theta(x_t, t)$$

Linearly combine denoisers from an unconditional distribution and a conditional distribution

# GANs: Learning to play a two-party game



Generative  
Adversarial  
Networks  
(GANs)

# Generative Adversarial Networks

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

Problem: Want to sample from complex, high-dimensional training distribution. No direct way to do this!

Solution: Sample from a simple distribution we can easily sample from, e.g. random noise. Learn transformation to training distribution.

# Generative Adversarial Networks

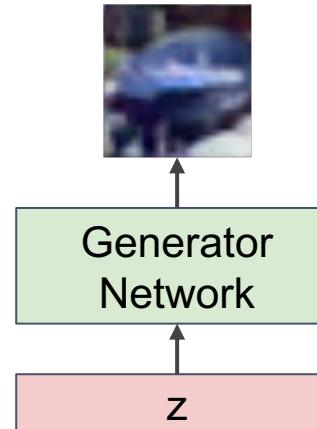
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Output: Sample from  
training distribution

Input: Random noise



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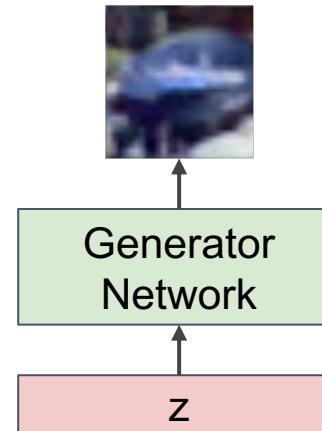
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Solution: Sample from a simple distribution we can easily sample from, e.g. random noise. Learn transformation to training distribution.

But we don't know which sample  $z$  maps to which training image -> can't learn by reconstructing training images

Output: Sample from training distribution

Input: Random noise



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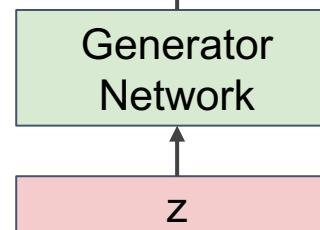
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Input: Random noise



Objective: generated images should look “real”

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Problem: Want to sample from complex, high-dimensional training distribution. No direct way to do this!

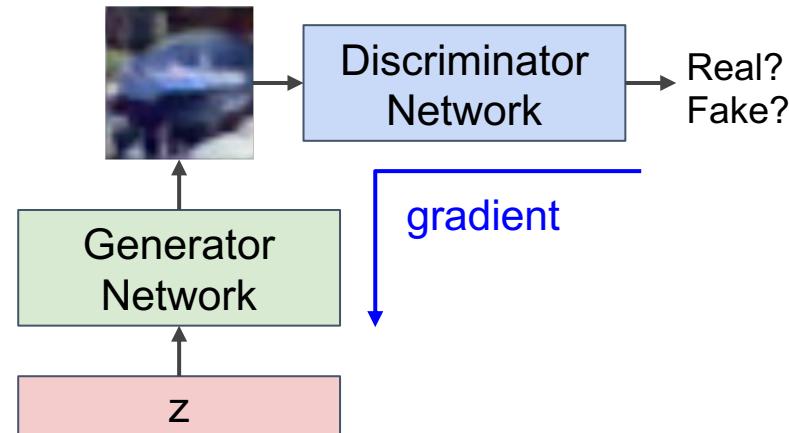
Solution: Sample from a simple distribution we can easily sample from, e.g. random noise. Learn transformation to training distribution.

But we don't know which sample  $z$  maps to which training image -> can't learn by reconstructing training images

Solution: Use a discriminator network to tell whether the generate image is within data distribution (“real”) or not

Output: Sample from training distribution

Input: Random noise



# Training GANs: Two-player game

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

**Discriminator network:** try to distinguish between real and fake images

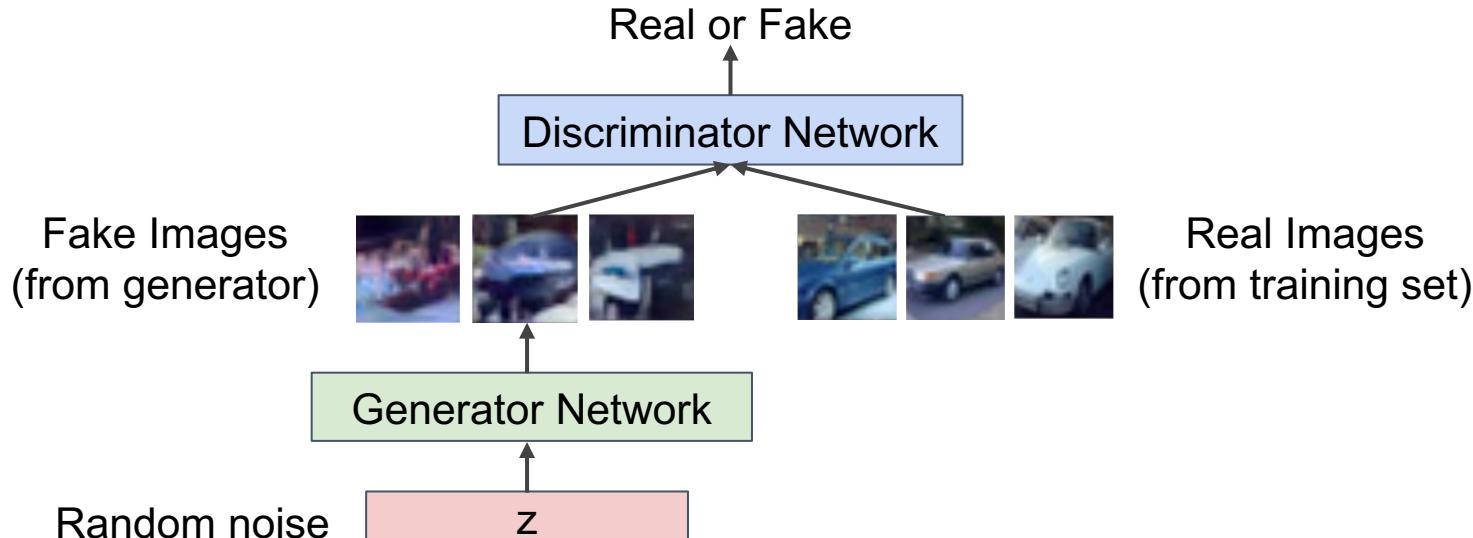
**Generator network:** try to fool the discriminator by generating real-looking images

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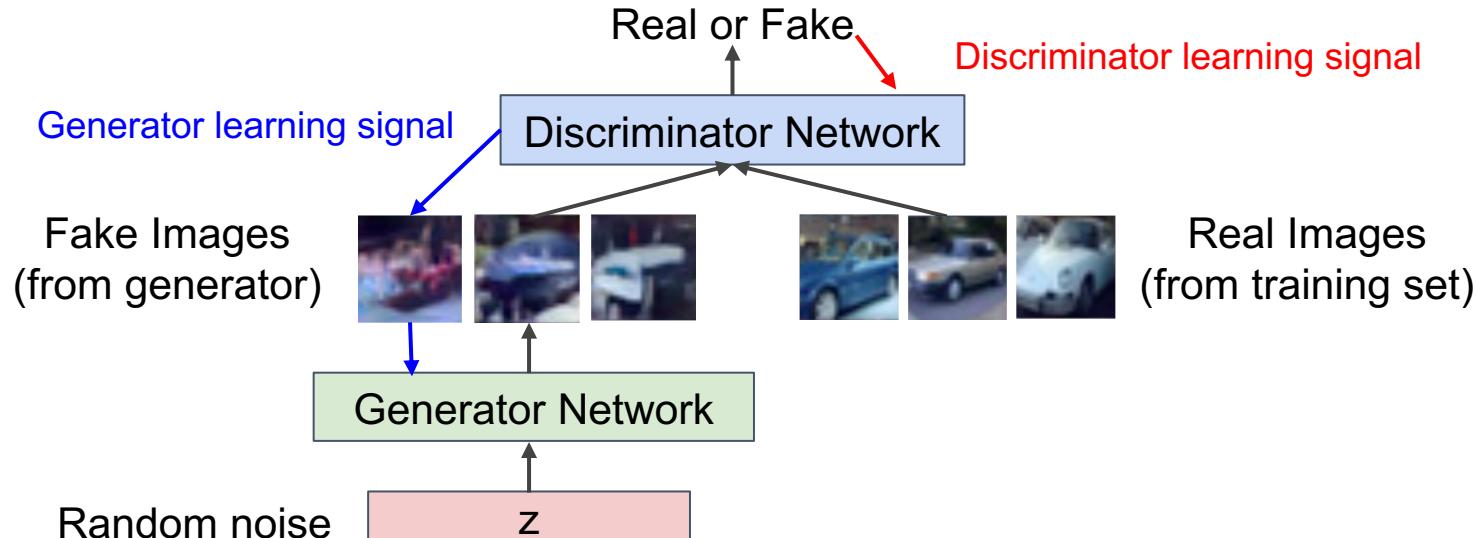
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# Training GANs: Two-player game

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

**Discriminator network:** try to distinguish between real and fake images

**Generator network:** try to fool the discriminator by generating real-looking images

Train jointly in **minimax game**

Minimax objective function:

$$\min_{\theta_g} \max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

Generator objective  
Discriminator objective

# Training GANs: Two-player game

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$$\min_{\theta_g} \max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log \underbrace{D_{\theta_d}(x)}_{\text{Discriminator output for real data } x} + \mathbb{E}_{z \sim p(z)} \log \underbrace{(1 - D_{\theta_d}(G_{\theta_g}(z)))}_{\text{Discriminator output for generated fake data } G(z)} \right]$$

Classify all real images as real

Classify all generated images as fake

Discriminator outputs likelihood in (0,1) of real image

# Training GANs: Two-player game

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Generator: learn to fool  
discriminator. Minimize  
 $\log(1 - p_{\theta_d}(x_{gen}))$

# Training GANs: Two-player game

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Train jointly in **minimax game**

Discriminator outputs likelihood in (0,1) of real image

Minimax objective function:

$$\min_{\theta_g} \max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

- Discriminator ( $\theta_d$ ) wants to **maximize objective** such that  $D(x)$  is close to 1 (real) and  $D(G(z))$  is close to 0 (fake)
- Generator ( $\theta_g$ ) wants to **minimize objective** such that  $D(G(z))$  is close to 1 (discriminator is fooled into thinking generated  $G(z)$  is real)

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Minimax objective function:

$$\min_{\theta_g} \max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

Alternate between:

1. **Gradient ascent** on discriminator

$$\max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

2. **Gradient descent** on generator

$$\min_{\theta_g} \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$

# Training GANs: Two-player game

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

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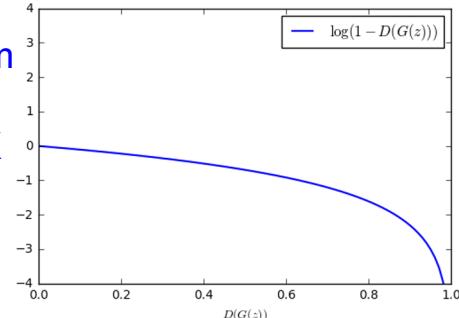
$$\max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

2. **Gradient descent** on generator

$$\min_{\theta_g} \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$

In practice, optimizing this generator objective does not work well!

When sample is likely fake, want to learn from it to improve generator (move to the right on X axis).



# Training GANs: Two-player game

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

Minimax objective function:

$$\min_{\theta_g} \max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

Alternate between:

1. **Gradient ascent** on discriminator

$$\max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

Gradient signal dominated by region where sample is already good

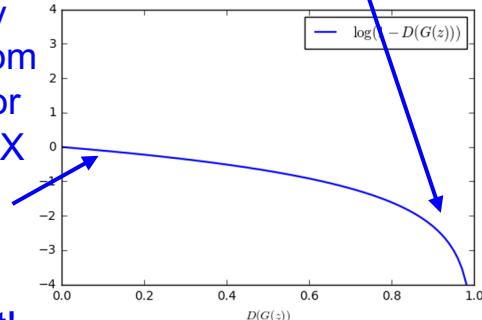
2. **Gradient descent** on generator

$$\min_{\theta_g} \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z)))$$

In practice, optimizing this generator objective does not work well!

When sample is likely fake, want to learn from it to improve generator (move to the right on X axis).

But gradient in this region is relatively flat!



# Training GANs: Two-player game

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

Minimax objective function:

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Alternate between:

1. **Gradient ascent** on discriminator

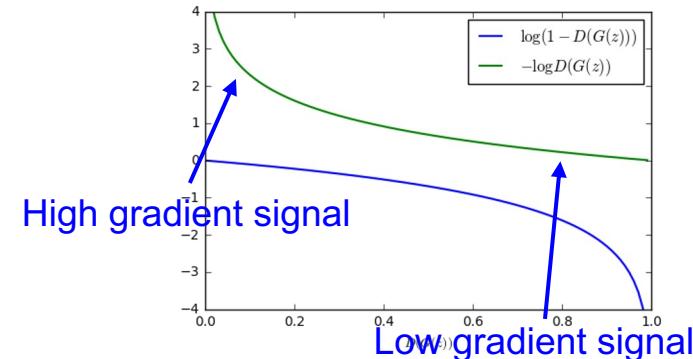
$$\max_{\theta_d} \left[ \mathbb{E}_{x \sim p_{data}} \log D_{\theta_d}(x) + \mathbb{E}_{z \sim p(z)} \log(1 - D_{\theta_d}(G_{\theta_g}(z))) \right]$$

2. Instead: **Gradient ascent** on generator, different objective

$$\max_{\theta_g} \mathbb{E}_{z \sim p(z)} \log(D_{\theta_d}(G_{\theta_g}(z)))$$

Instead of minimizing likelihood of discriminator being correct, now maximize likelihood of discriminator being wrong.

Same objective of fooling discriminator, but now higher gradient signal for bad samples => works much better! Standard in practice.



# Training GANs: Two-player game

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

## Putting it together: GAN training algorithm

**for** number of training iterations **do**

**for**  $k$  steps **do**

- Sample minibatch of  $m$  noise samples  $\{\mathbf{z}^{(1)}, \dots, \mathbf{z}^{(m)}\}$  from noise prior  $p_g(\mathbf{z})$ .
- Sample minibatch of  $m$  examples  $\{\mathbf{x}^{(1)}, \dots, \mathbf{x}^{(m)}\}$  from data generating distribution  $p_{\text{data}}(\mathbf{x})$ .
- Update the discriminator by ascending its stochastic gradient:

$$\nabla_{\theta_d} \frac{1}{m} \sum_{i=1}^m \left[ \log D_{\theta_d}(\mathbf{x}^{(i)}) + \log(1 - D_{\theta_d}(G_{\theta_g}(\mathbf{z}^{(i)}))) \right]$$

**end for**

- Sample minibatch of  $m$  noise samples  $\{\mathbf{z}^{(1)}, \dots, \mathbf{z}^{(m)}\}$  from noise prior  $p_g(\mathbf{z})$ .
- Update the generator by ascending its stochastic gradient (improved objective):

$$\nabla_{\theta_g} \frac{1}{m} \sum_{i=1}^m \log(D_{\theta_d}(G_{\theta_g}(\mathbf{z}^{(i)})))$$

**end for**

# Training GANs: Two-player game

Ian Goodfellow et al., "Generative Adversarial Nets", NIPS 2014

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- Update the generator by ascending its stochastic gradient (improved objective):

$$\nabla_{\theta_g} \frac{1}{m} \sum_{i=1}^m \log(D_{\theta_d}(G_{\theta_g}(z^{(i)})))$$

**end for**

Some find  $k=1$  more stable, others use  $k > 1$ , no best rule.

Followup work (e.g. Wasserstein GAN, BEGAN) alleviates this problem, better stability!

Update discriminator

Update generator

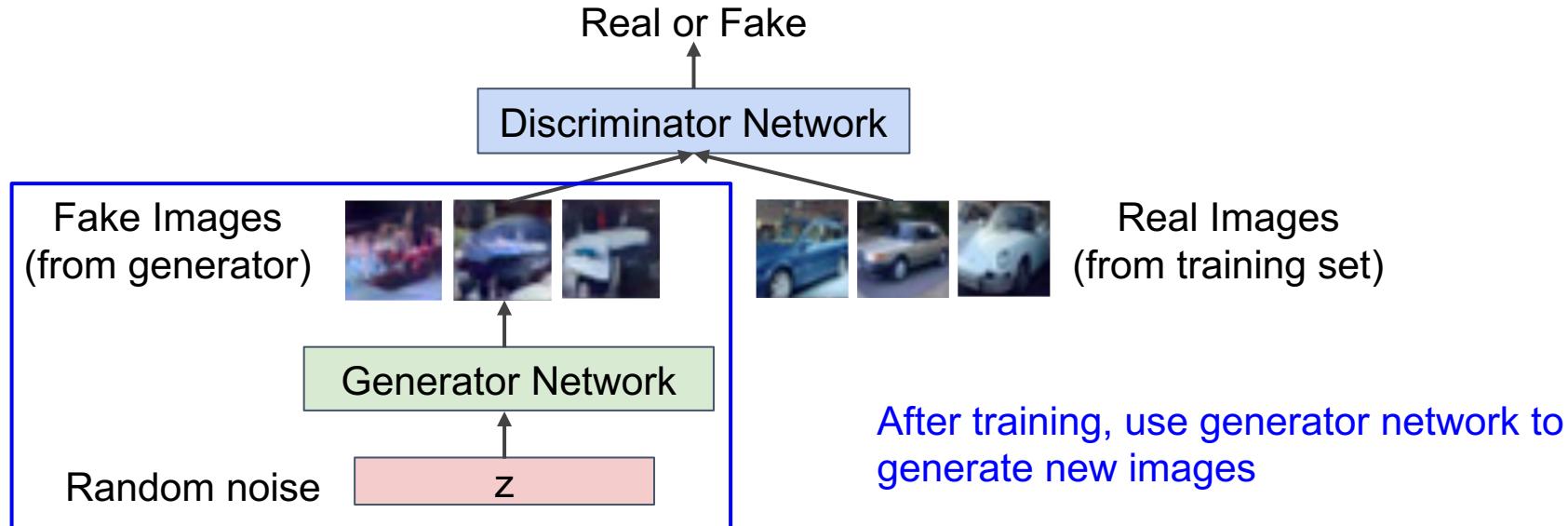
Arjovsky et al. "Wasserstein gan." arXiv preprint arXiv:1701.07875 (2017)

Berthelot, et al. "Began: Boundary equilibrium generative adversarial networks." arXiv preprint arXiv:1703.10717 (2017)

# Training GANs: Two-player game

Ian Goodfellow et al., “Generative Adversarial Nets”, NIPS 2014

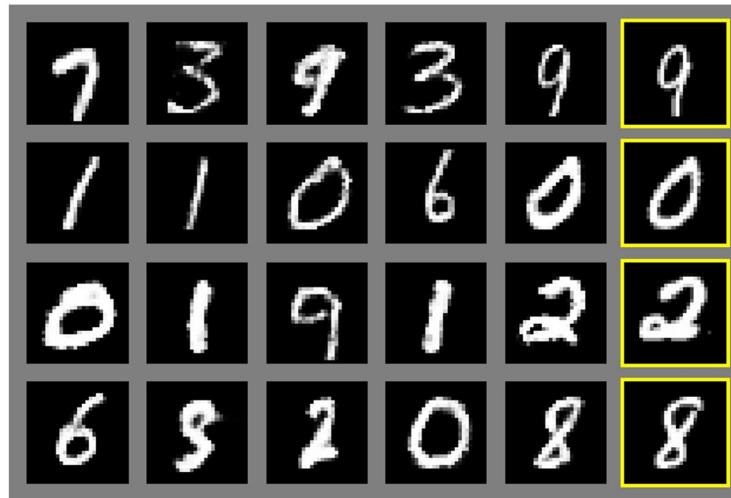
**Generator network:** try to fool the discriminator by generating real-looking images  
**Discriminator network:** try to distinguish between real and fake images



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# Generative Adversarial Nets

Generated samples



Nearest neighbor from training set

# Generative Adversarial Nets

Generated samples (CIFAR-10)



Nearest neighbor from training set

# Generative Adversarial Nets: Convolutional Architectures

Generator is an upsampling network with fractionally-strided convolutions  
Discriminator is a convolutional network

Architecture guidelines for stable Deep Convolutional GANs

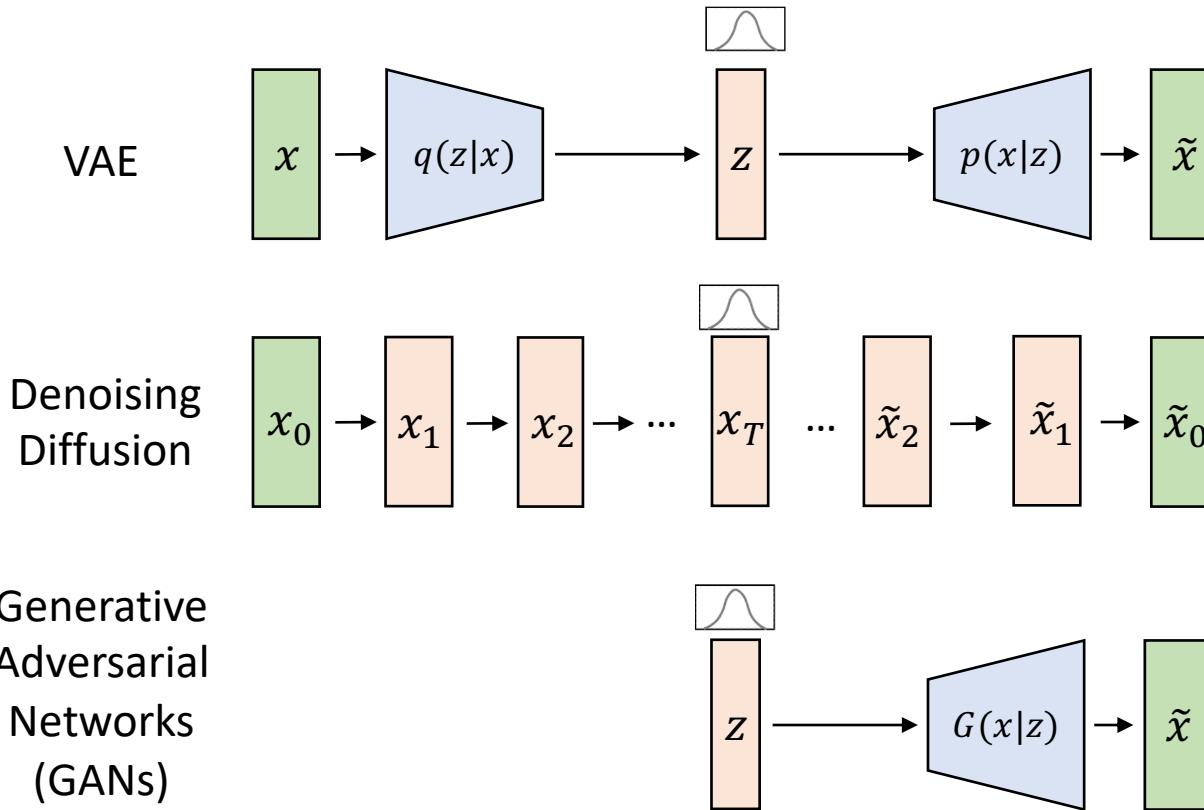
- Replace any pooling layers with strided convolutions (discriminator) and fractional-strided convolutions (generator).
- Use batchnorm in both the generator and the discriminator.
- Remove fully connected hidden layers for deeper architectures.
- Use ReLU activation in generator for all layers except for the output, which uses Tanh.
- Use LeakyReLU activation in the discriminator for all layers.

# 2019: BigGAN



Brock et al., 2019

# Deep Generative Models



# Generative Models: Closing Thoughts

- Learn without supervision = ability to leverage large, raw dataset
- Realism: Generate plausible samples given dataset
- Diversity: Generate diverse samples (avoid mode collapse)
- Controllability: Generate based on instruction / conditioning
- Healthy combination of theory and deep learning magic
- Generative Model is extremely hot in 2023. More will come ...

## Supervised Learning

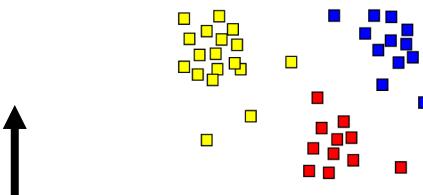
- Train Input:  $\{X, Y\}$
- Learning output:  
 $f : X \rightarrow Y, P(y|x)$
- e.g. classification



Sheep  
Dog  
Cat  
Lion  
Giraffe

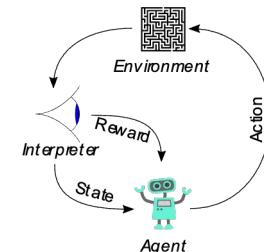
## Unsupervised Learning

- Input:  $\{X\}$
- Learning output:  $P(x)$
- Example: Clustering, density estimation, generative modeling



## Reinforcement Learning

- Evaluative feedback in the form of **reward**
- No supervision on the right action

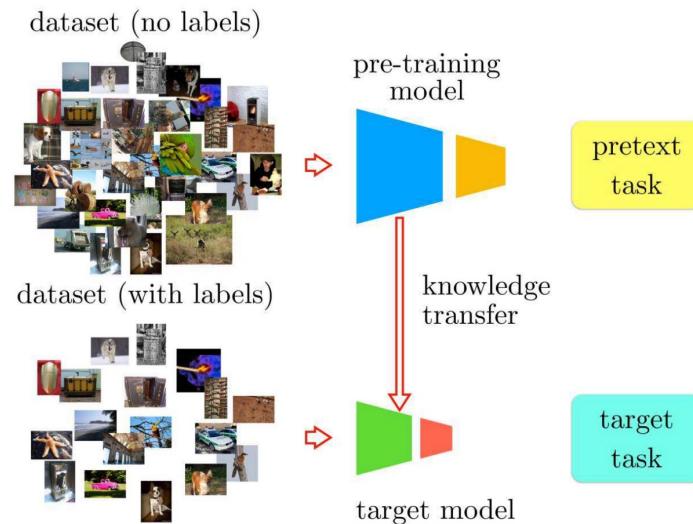


**Self-Supervised Learning:**  
Create your own supervision

# Self-supervised Learning

In short: still supervised learning, with two important distinctions:

1. Learn from labels generated *autonomously* instead of human annotations.
2. The goal is to learn *good representations* for *other target tasks*.



# Self-supervised pretext tasks

Example: learn to predict image transformations / complete corrupted images

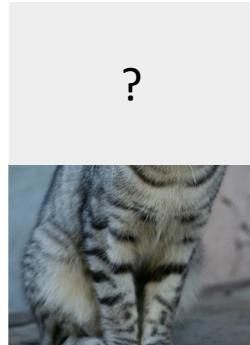
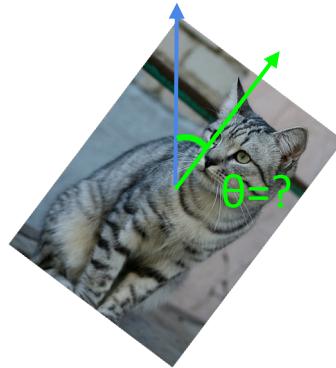


image completion



rotation prediction



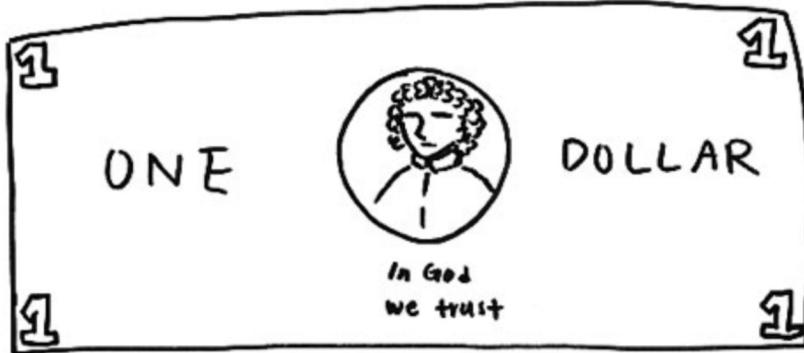
“jigsaw puzzle”



colorization

1. Solving the pretext tasks allow the model to learn good features.
2. We can automatically generate labels for the pretext tasks.

# Generative vs. Self-supervised Learning



Left: Drawing of a dollar bill from memory. Right: Drawing subsequently made with a dollar bill present. Image source: [Epstein, 2016](#)

Learning to generate pixel-level details is often unnecessary; learn high-level semantic features with pretext tasks instead

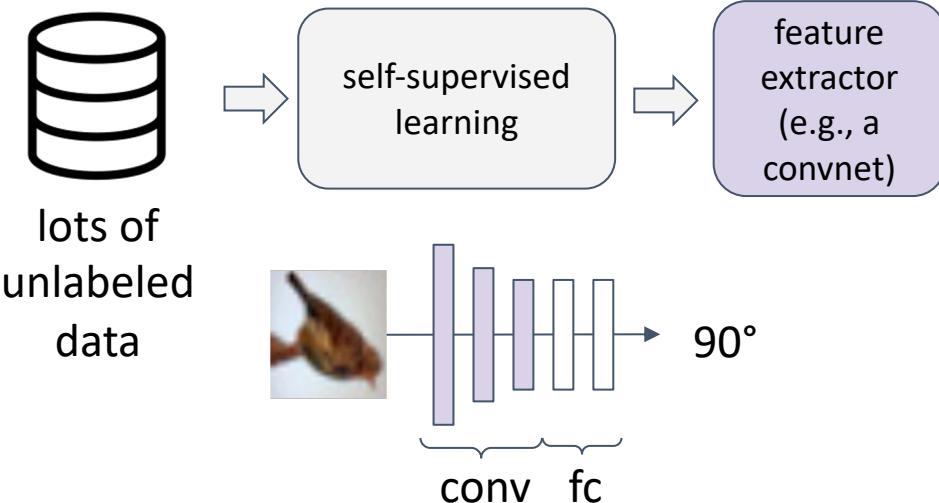
Source: [Anand, 2020](#)

# How to evaluate a self-supervised learning method?

We usually don't care about the performance of the self-supervised learning task, e.g., we don't care if the model learns to predict image rotation perfectly.

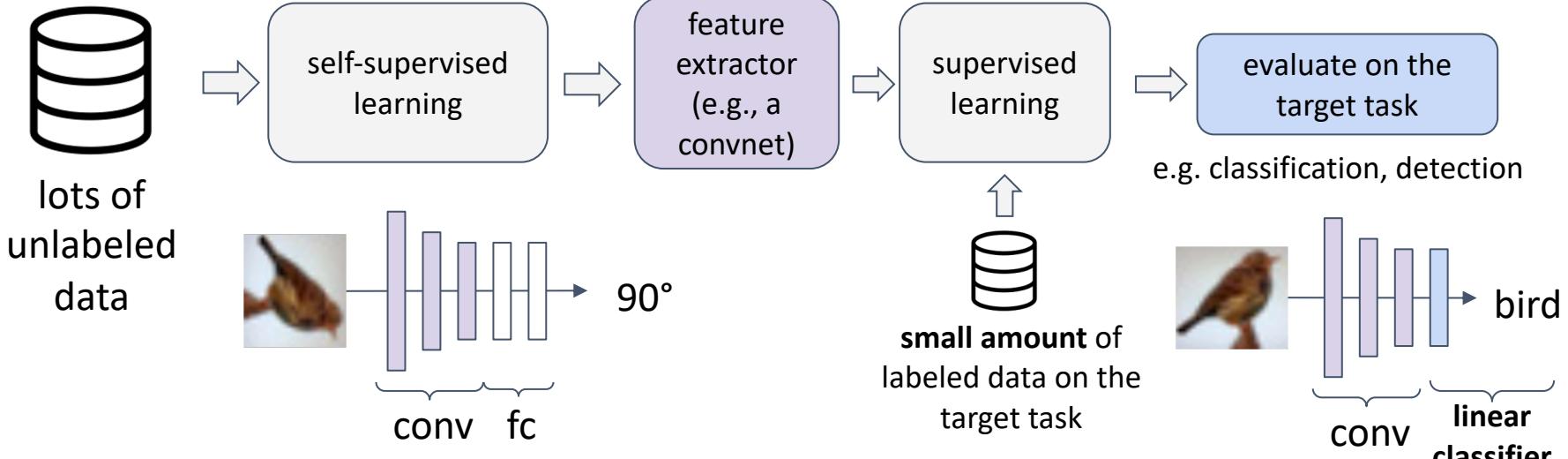
Evaluate the learned feature encoders on downstream *target tasks*

# How to evaluate a self-supervised learning method?



1. Learn good feature extractors from self-supervised pretext tasks, e.g., predicting image rotations

# How to evaluate a self-supervised learning method?



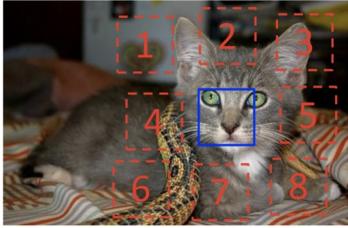
1. Learn good feature extractors from self-supervised pretext tasks, e.g., predicting image rotations

2. Attach a shallow network on the feature extractor; train the shallow network on the target task with small amount of labeled data

# Broader picture

## Today's lecture

### computer vision



Doersch et al., 2015

### language modeling

#### Language Models are Few-Shot Learners

Tom B. Brown\* Benjamin Mann\* Nick Ryder\* Melanie Subbiah\*  
Jared Kaplan† Prafulla Dhariwal Arvind Neelakantan Pranav Shyam Girish Sastry  
Amanda Askell Sandhini Agarwal Ariel Herbert-Voss Gretchen Krueger Tom Henighan  
Rewon Child Aditya Ramesh Daniel M. Ziegler Jeffrey Wu Clemens Winter  
Christopher Hesse Mark Chen Eric Sigler Mateusz Litwin Scott Gray  
Benjamin Chess Jack Clark Christopher Berner  
Sam McCandlish Alec Radford Ilya Sutskever Dario Amodei  
OpenAI

#### Abstract

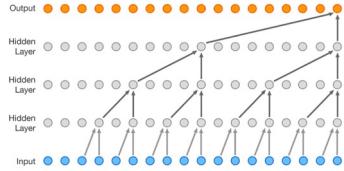
Recent work has demonstrated substantial gains on many NLP tasks and benchmarks by pre-training on a large corpus of text followed by fine-tuning on a specific task. While typically task-agnostic in architecture, this method still requires task-specific fine-tuning datasets of thousands or tens of thousands of examples. By contrast, humans can generally perform a new language task from only a few examples. From this perspective, scaling word embeddings with context (‘‘synonyms’’) appears to struggle to do. Here we show that scaling of language models greatly improves task-agnostic few-shot performance, sometimes even reaching competitiveness with prior state-of-the-art fine-tuning approaches. Specifically, we train GPT-3, an autoregressive language model with 175 billion parameters, 10x more than any previous non-sparse language model, and test its performance in the few-shot setting. Furthermore, GPT-3 is applied without task-specific fine-tuning or fine-tuning, with tasks and few-shot demonstrations specified purely via text interaction with the model. GPT-3 achieves strong performance on many NLP datasets, including translation, question answering, and cloze tasks, as well as several tasks that require on-the-fly reasoning or domain adaptation, such as unscrambling words, using a novel word in a sentence, or performing 3-digit arithmetic. At the same time, we find identical requirements where GPT-3 fails to learn quickly, as well as some datasets where GPT-3 faces methodological issues related to training on large web corpora. Finally, we find that GPT-3 can generate samples of news articles which human evaluators have difficulty distinguishing from articles written by humans. We discuss broader societal impacts of this finding and of GPT-3 in general.

### robot / reinforcement learning



Dense Object Net (Florence and Manuelli et al., 2018)

### speech synthesis



Wavenet (van den Oord et al., 2016)

• • •

# Today's Agenda

## **Pretext tasks from image transformations**

- Rotation, inpainting, rearrangement, coloring

## **Contrastive representation learning**

- Intuition and formulation
- Instance contrastive learning: SimCLR and MOCO
- Sequence contrastive learning: CPC

# Today's Agenda

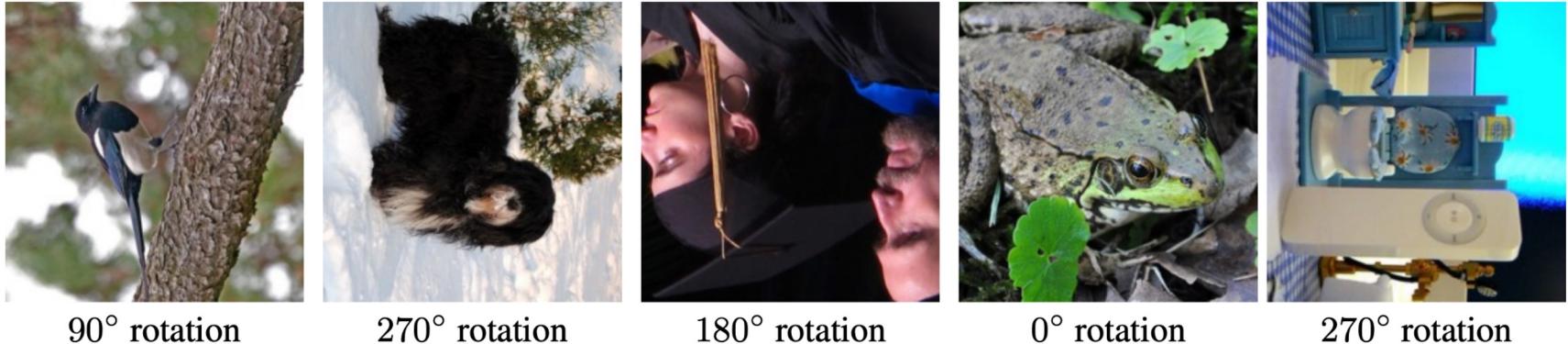
## **Pretext tasks from image transformations**

- Rotation, inpainting, rearrangement, coloring

## **Contrastive representation learning**

- Intuition and formulation
- Instance contrastive learning: SimCLR and MOCO
- Sequence contrastive learning: CPC

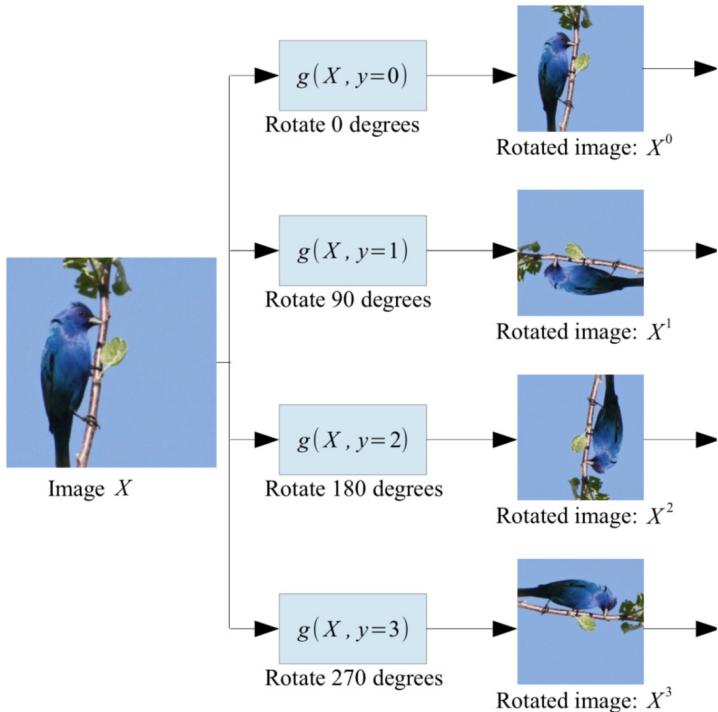
# Pretext task: predict rotations



**Hypothesis:** a model could recognize the correct rotation of an object only if it has the “visual commonsense” of what the object should look like unperturbed.

(Image source: [Gidaris et al. 2018](#))

# Pretext task: predict rotations

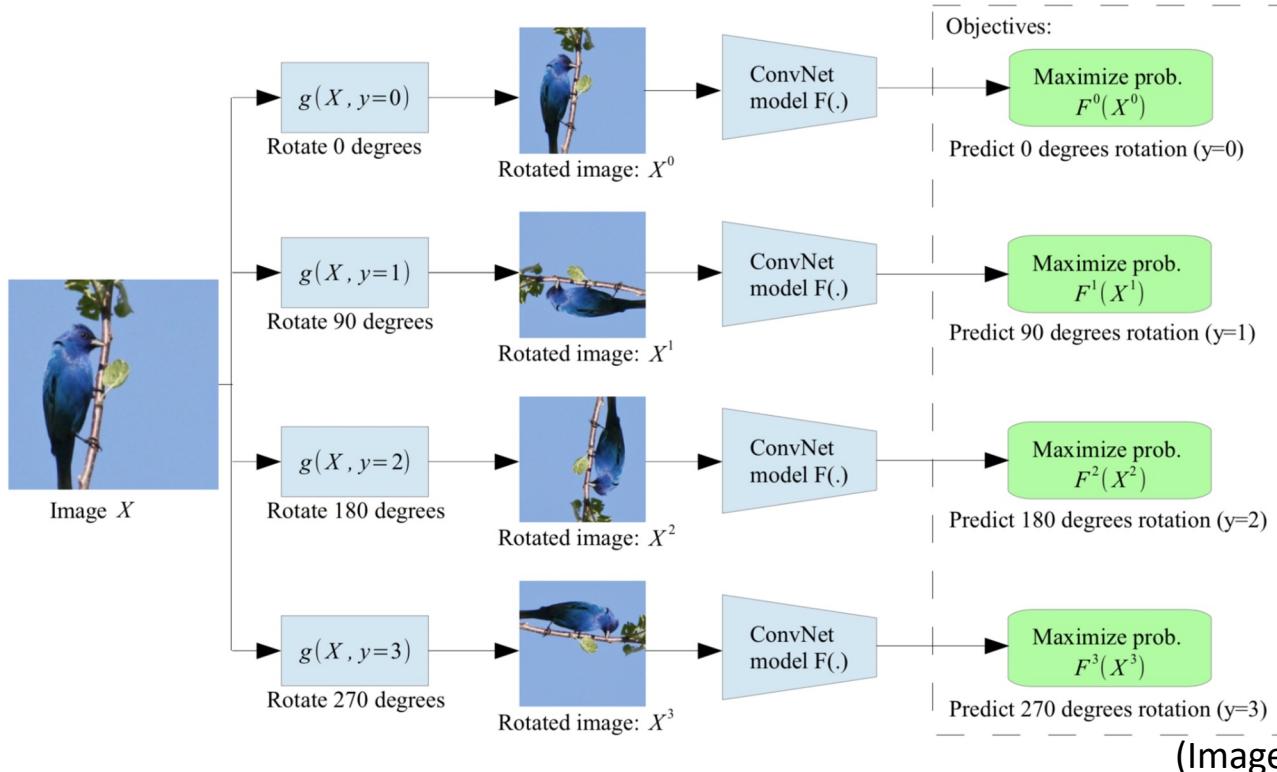


Self-supervised learning by rotating the entire input images.

The model learns to predict which rotation is applied (4-way classification)

(Image source: [Gidaris et al. 2018](#))

# Pretext task: predict rotations

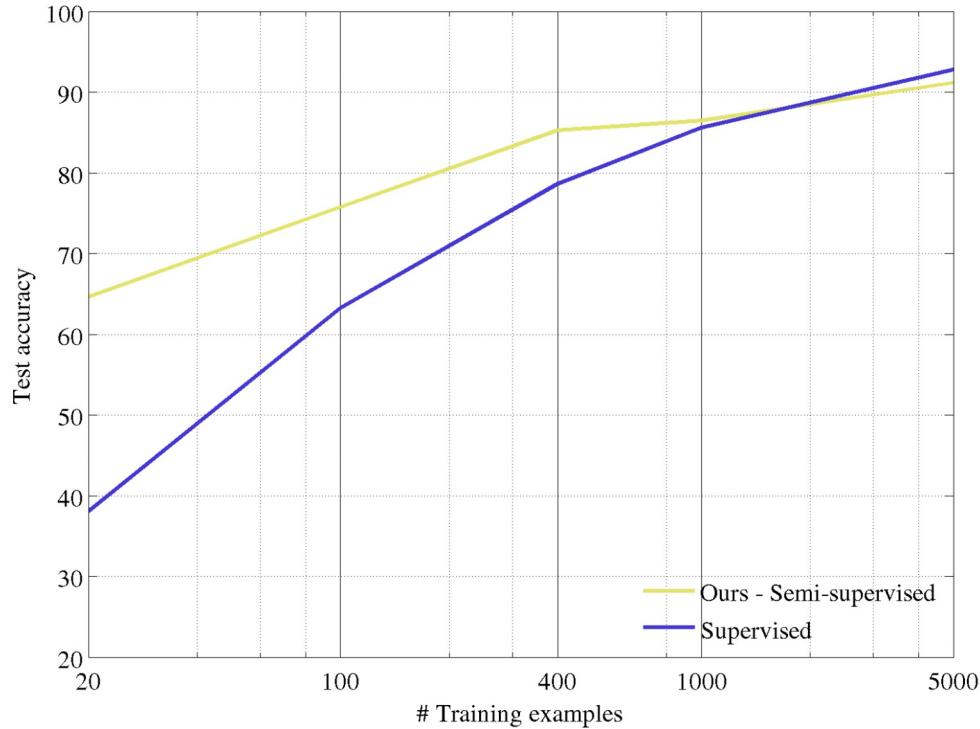


Self-supervised learning by rotating the entire input images.

The model learns to predict which rotation is applied (4-way classification)

(Image source: [Gidaris et al. 2018](#))

# Evaluation on semi-supervised learning



Self-supervised learning on  
**CIFAR10** (entire training set).

Freeze conv1 + conv2  
Learn **conv3 + linear** layers with  
subset of labeled CIFAR10 data  
(classification).

(Image source: [Gidaris et al. 2018](#))

# Transfer learned features to supervised learning

	Classification (%mAP)	Detection (%mAP)	Segmentation (%mIoU)
Trained layers	fc6-8	all	all
ImageNet labels	78.9	79.9	56.8
Random		53.3	43.4
Random rescaled Krähenbühl et al. (2015)	39.2	56.6	45.6
Egomotion (Agrawal et al., 2015)	31.0	54.2	43.9
Context Encoders (Pathak et al., 2016b)	34.6	56.5	44.5
Tracking (Wang & Gupta, 2015)	55.6	63.1	47.4
Context (Doersch et al., 2015)	55.1	65.3	51.1
Colorization (Zhang et al., 2016a)	61.5	65.6	46.9
BIGAN (Donahue et al., 2016)	52.3	60.1	46.9
Jigsaw Puzzles (Noroozi & Favaro, 2016)	-	67.6	53.2
NAT (Bojanowski & Joulin, 2017)	56.7	65.3	49.4
Split-Brain (Zhang et al., 2016b)	63.0	67.1	46.7
ColorProxy (Larsson et al., 2017)		65.9	
Counting (Noroozi et al., 2017)	-	67.7	51.4
(Ours) RotNet	<b>70.87</b>	<b>72.97</b>	<b>54.4</b>
			<b>39.1</b>

Pretrained with full  
ImageNet supervision

No pretraining

Self-supervised learning on  
**ImageNet** (entire training  
set) with AlexNet.

Finetune on labeled data  
from **Pascal VOC 2007**.

Self-supervised learning with rotation  
prediction

source: [Gidaris et al. 2018](#)

# Visualize learned visual attentions



Conv1  $27 \times 27$    Conv3  $13 \times 13$    Conv5  $6 \times 6$

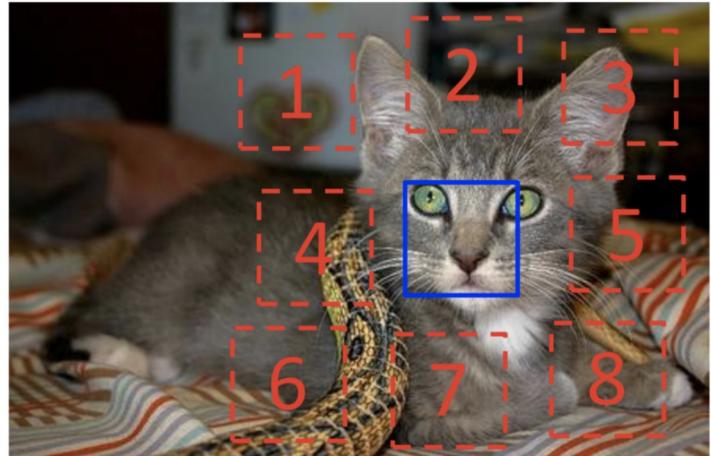
Conv1  $27 \times 27$    Conv3  $13 \times 13$    Conv5  $6 \times 6$

(a) Attention maps of supervised model

(b) Attention maps of our self-supervised model

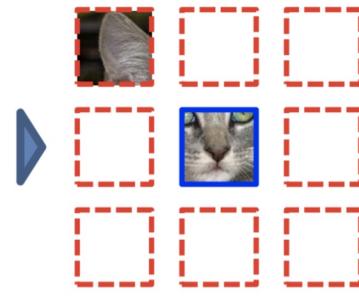
(Image source: [Gidaris et al. 2018](#))

# Pretext task: predict relative patch locations



$$X = (\text{cat eye}, \text{ear}) ; Y = 3$$

Example:



Question 1:

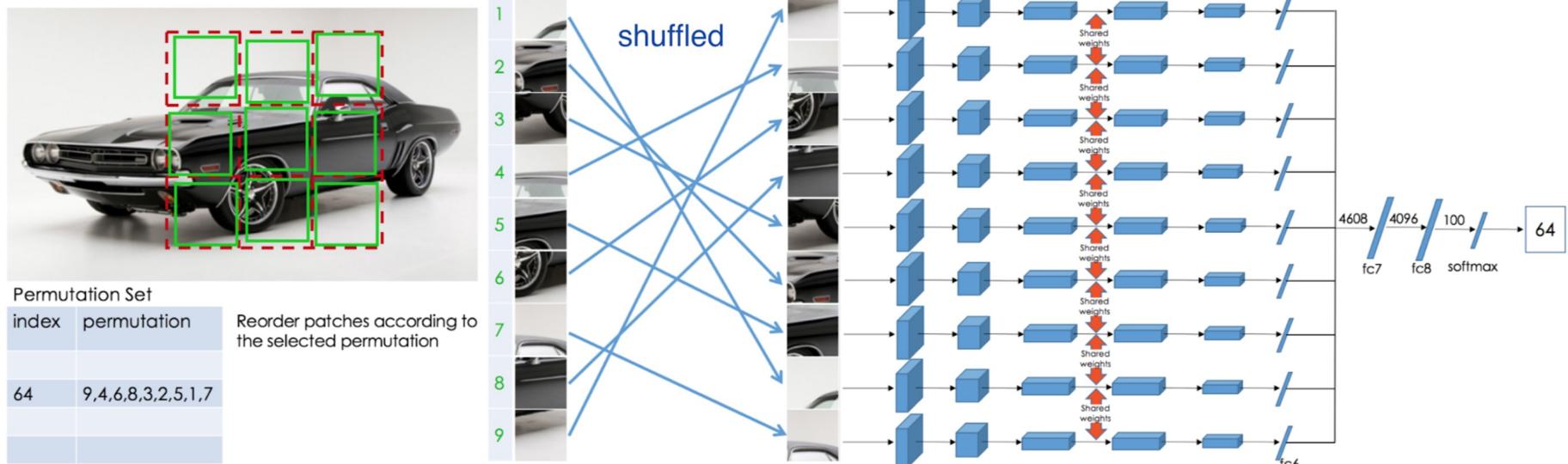


Question 2:



(Image source: [Doersch et al., 2015](#))

# Pretext task: solving “jigsaw puzzles”



(Image source: [Noroozi & Favaro, 2016](#))

# Transfer learned features to supervised learning

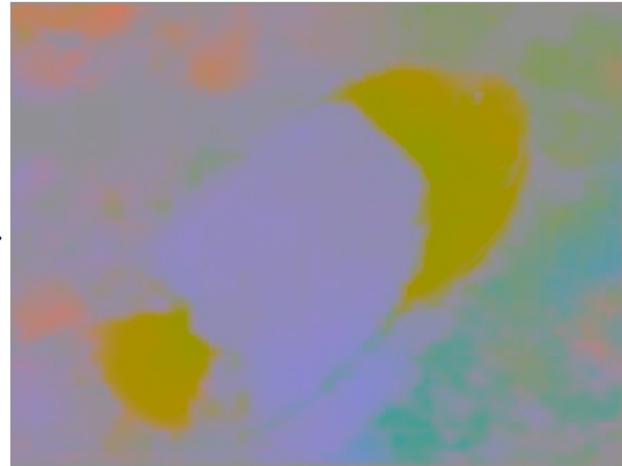
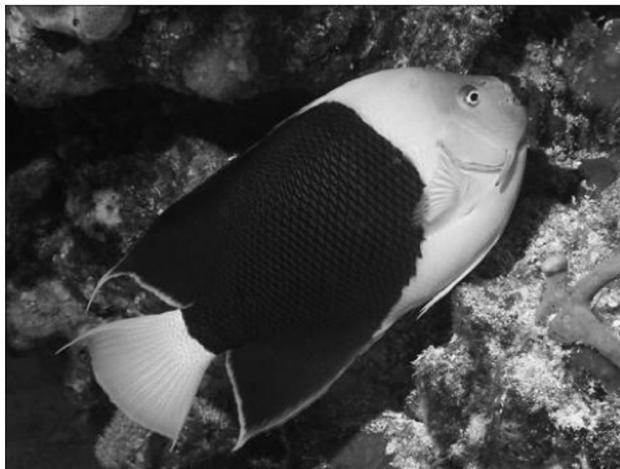
Table 1: Results on PASCAL VOC 2007 Detection and Classification. The results of the other methods are taken from Pathak *et al.* [30].

Method	Pretraining time	Supervision	Classification	Detection	Segmentation
Krizhevsky <i>et al.</i> [25]	3 days	1000 class labels	<b>78.2%</b>	<b>56.8%</b>	<b>48.0%</b>
Wang and Gupta[39]	1 week	motion	58.4%	44.0%	-
Doersch <i>et al.</i> [10]	4 weeks	context	55.3%	46.6%	-
Pathak <i>et al.</i> [30]	14 hours	context	56.5%	44.5%	29.7%
Ours	2.5 days	context	<b>67.6%</b>	<b>53.2%</b>	<b>37.6%</b>

“Ours” is feature learned from solving image Jigsaw puzzles (Noroozi & Favaro, 2016). Doersch et al. is the method with relative patch location

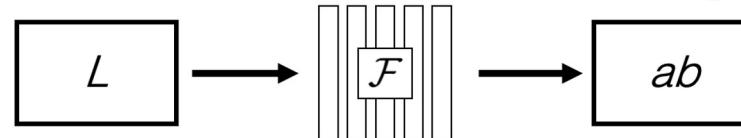
(source: [Noroozi & Favaro, 2016](#))

# Pretext task: image coloring



Grayscale image:  $L$  channel

$$\mathbf{X} \in \mathbb{R}^{H \times W \times 1}$$



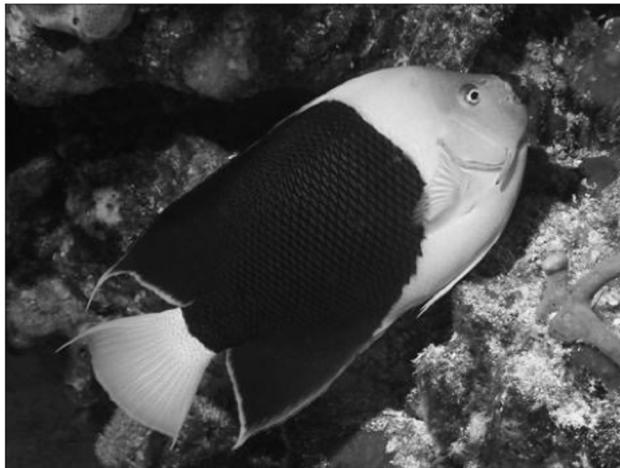
Color information:  $ab$  channels

$$\hat{\mathbf{Y}} \in \mathbb{R}^{H \times W \times 2}$$

5

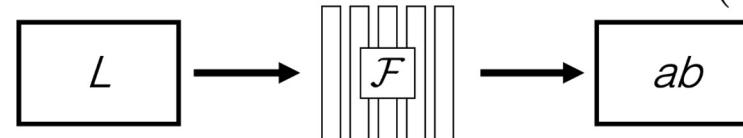
Source: Richard Zhang / Phillip Isola

# Pretext task: image coloring



Grayscale image:  $L$  channel

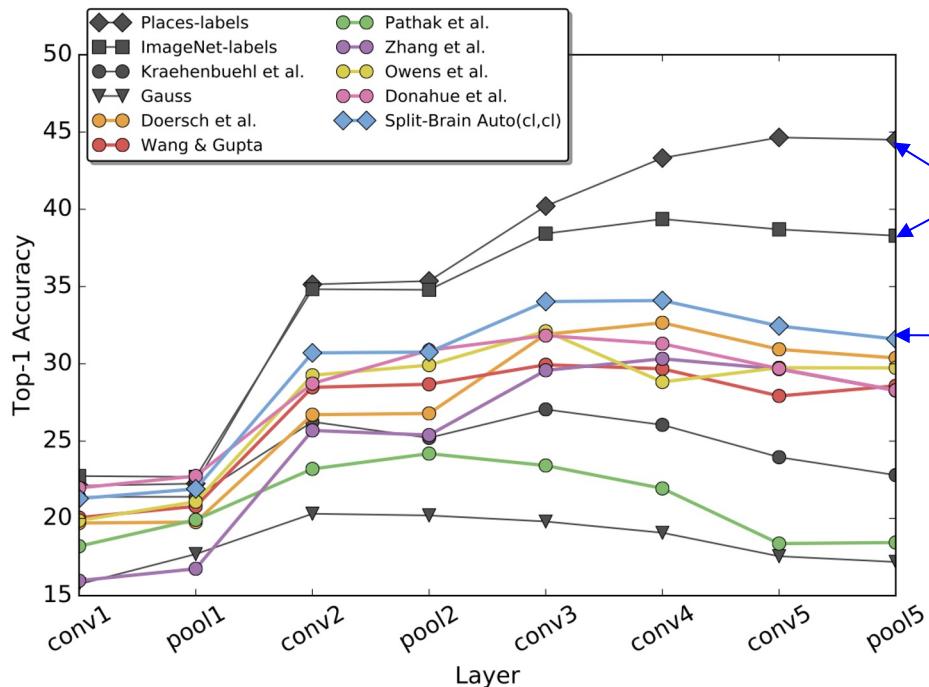
$$\mathbf{X} \in \mathbb{R}^{H \times W \times 1}$$



Concatenate  $(L, ab)$  channels  
 $(\mathbf{X}, \widehat{\mathbf{Y}})$

Source: Richard Zhang / Phillip Isola

# Transfer learned features to supervised learning



supervised

this paper

Self-supervised learning on  
**ImageNet** (entire training set).

Use *concatenated features*  
from  $F_1$  and  $F_2$

Labeled data is from the  
**Places** (Zhou 2016).

Source: [Zhang et al., 2017](#)

# Pretext task: image coloring



Source: Richard Zhang / Phillip Isola

# Pretext task: image coloring



Source: Richard Zhang / Phillip Isola

# Pretext task: video coloring

Idea: model the *temporal coherence* of colors in videos

reference frame



$t = 0$

how should I color these frames?



$t = 1$



$t = 2$



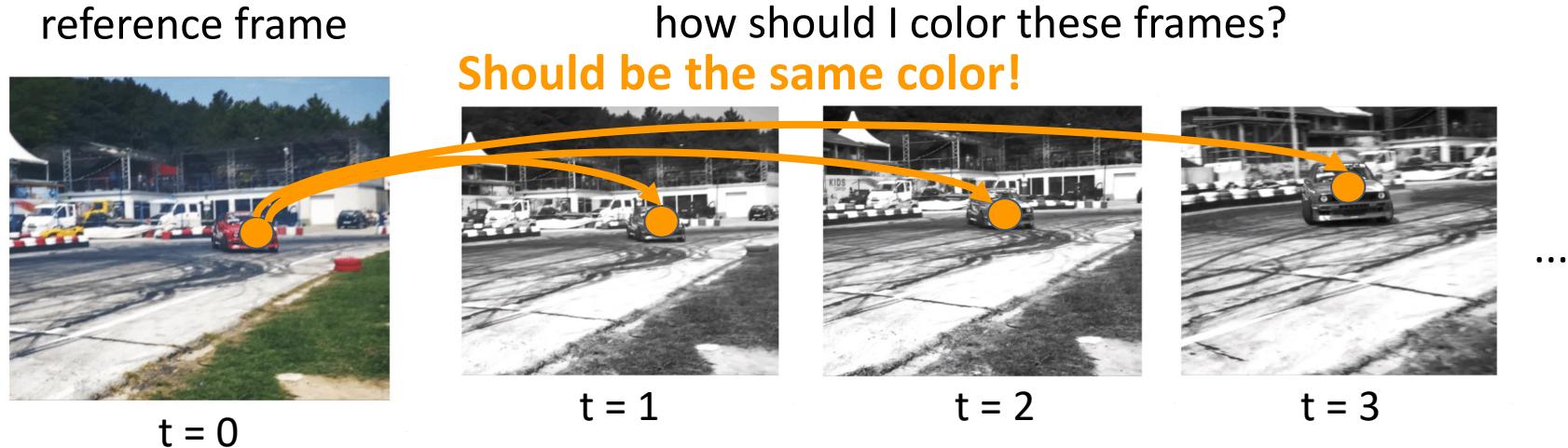
$t = 3$

...

Source: [Vondrick et al., 2018](#)

# Pretext task: video coloring

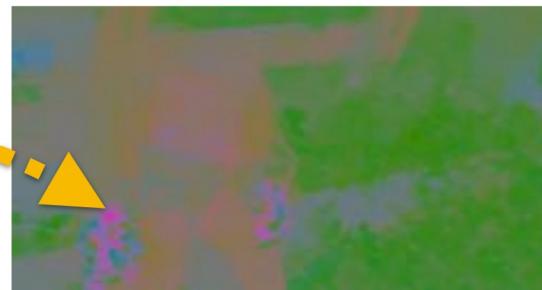
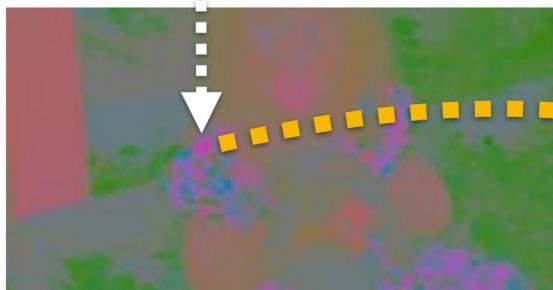
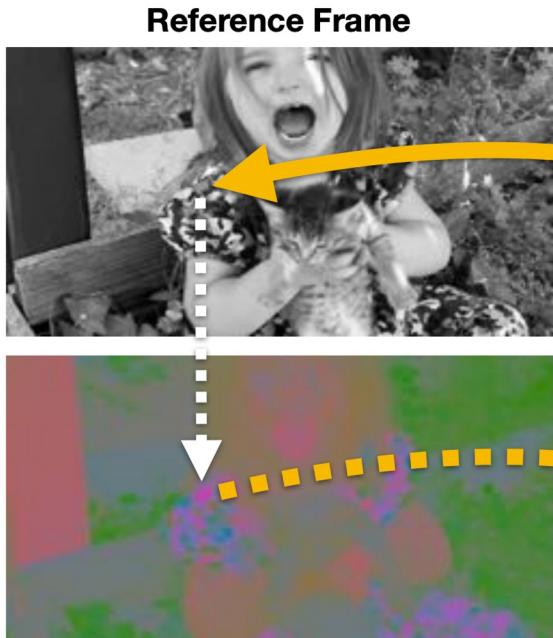
Idea: model the *temporal coherence* of colors in videos



**Hypothesis:** learning to color video frames should allow model to learn to track regions or objects without labels!

Source: [Vondrick et al., 2018](#)

# Learning to color videos



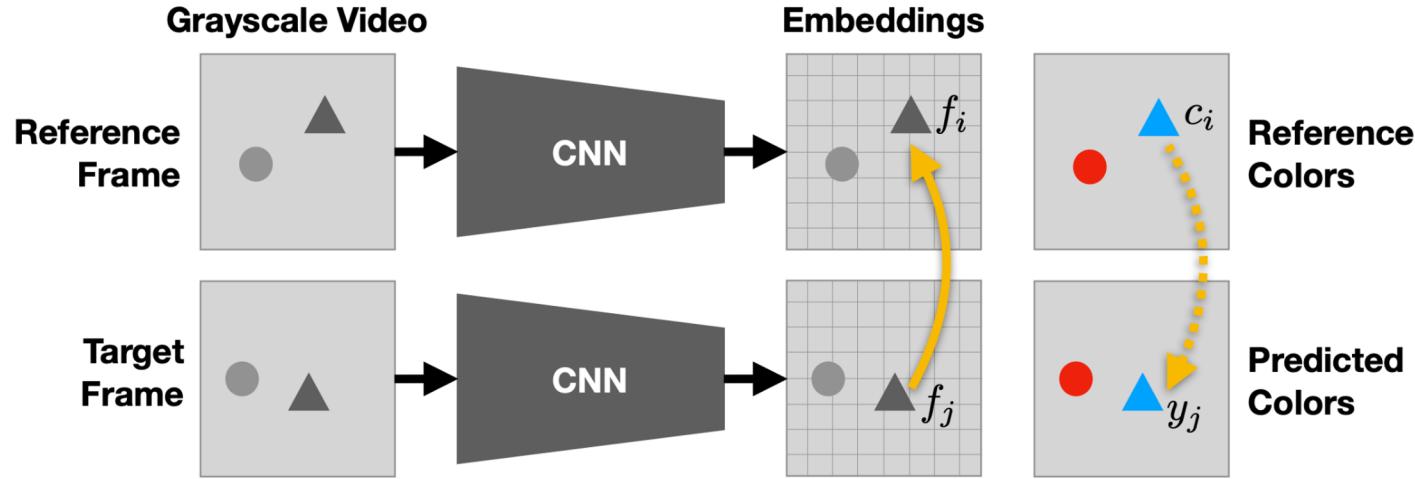
**Learning objective:**

Establish mappings between reference and target frames in a learned feature space.

Use the mapping as “pointers” to copy the correct color (LAB).

Source: [Vondrick et al., 2018](#)

# Learning to color videos

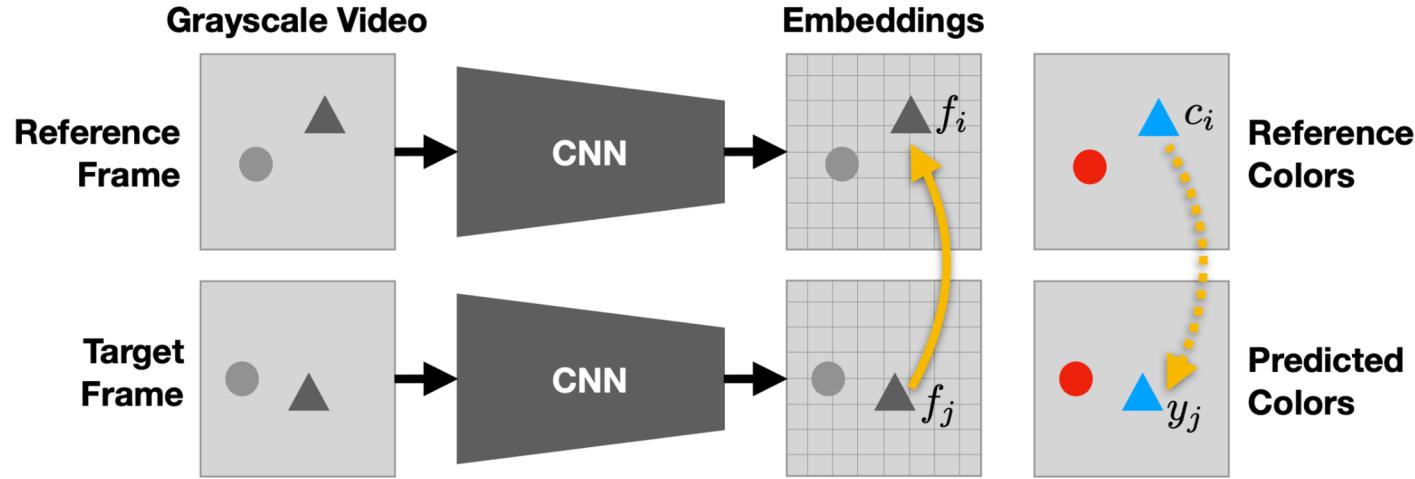


attention map on the reference frame

$$A_{ij} = \frac{\exp(f_i^T f_j)}{\sum_k \exp(f_k^T f_j)}$$

Source: [Vondrick et al., 2018](#)

# Learning to color videos



attention map on the reference frame

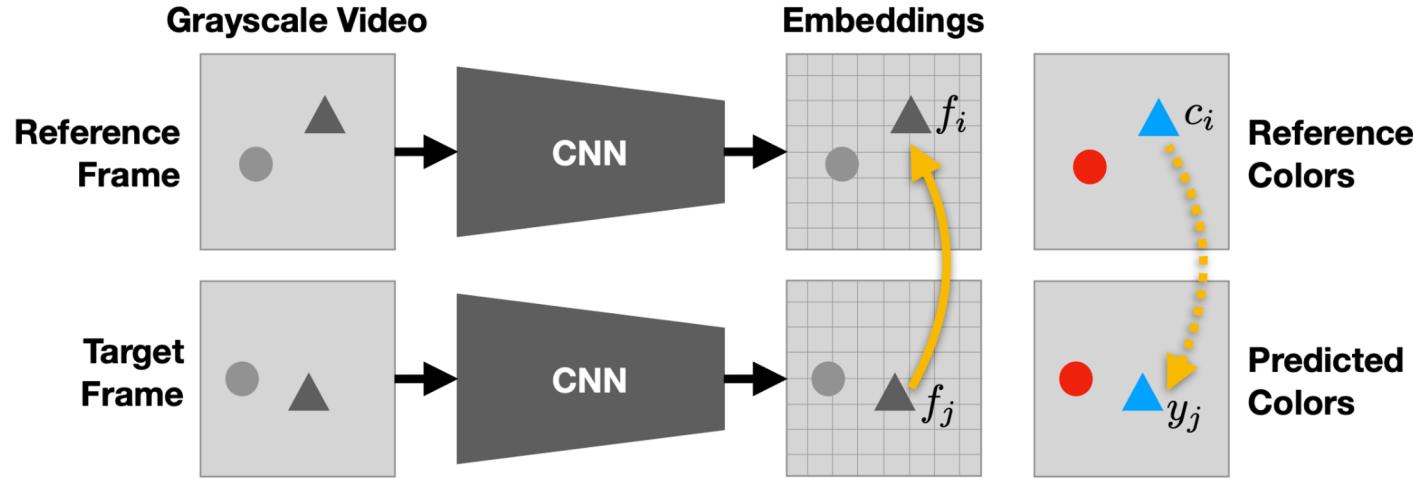
predicted color = weighted sum of the reference color

$$A_{ij} = \frac{\exp(f_i^T f_j)}{\sum_k \exp(f_k^T f_j)}$$

$$y_j = \sum_i A_{ij} c_i$$

Source: [Vondrick et al., 2018](#)

# Learning to color videos



attention map on the reference frame

$$A_{ij} = \frac{\exp(f_i^T f_j)}{\sum_k \exp(f_k^T f_j)}$$

predicted color = weighted sum of the reference color

$$y_j = \sum_i A_{ij} c_i$$

loss between predicted color and ground truth color

$$\min_{\theta} \sum_j \mathcal{L}(y_j, c_j)$$

Source: [Vondrick et al., 2018](#)

# Colorizing videos (qualitative)

reference frame



target frames (gray)



predicted color



Source: [Google AI blog post](#)

# Colorizing videos (qualitative)

reference frame



target frames (gray)



predicted color



Source: [Google AI blog post](#)

# Tracking emerges from colorization

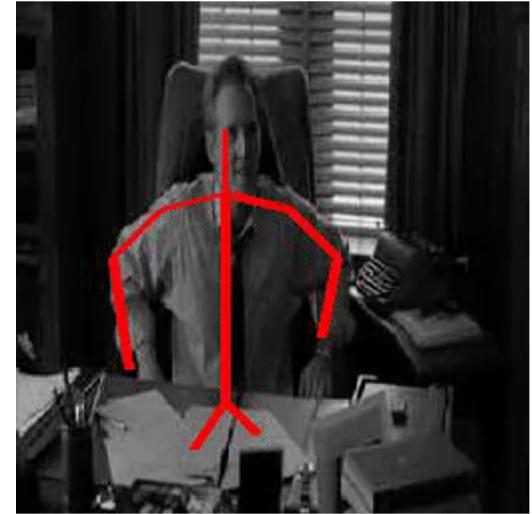
Propagate segmentation masks using learned attention



Source: [Google AI blog post](#)

# Tracking emerges from colorization

Propagate pose keypoints using learned attention



Source: [Google AI blog post](#)

# Summary: pretext tasks from image transformations

- Pretext tasks focus on “visual common sense”, e.g., predict rotations, inpainting, rearrangement, and colorization.
- The models are forced learn good features about natural images, e.g., semantic representation of an object category, in order to solve the pretext tasks.
- We don’t care about the performance of these pretext tasks, but rather how useful the learned features are for downstream tasks (classification, detection, segmentation).

# Summary: pretext tasks from image transformations

- Pretext tasks focus on “visual common sense”, e.g., predict rotations, inpainting, rearrangement, and colorization.
- The models are forced learn good features about natural images, e.g., semantic representation of an object category, in order to solve the pretext tasks.
- We don’t care about the performance of these pretext tasks, but rather how useful the learned features are for downstream tasks (classification, detection, segmentation).
- Problems: 1) coming up with individual pretext tasks is tedious, and 2) the learned representations may not be general.

# Pretext tasks from image transformations

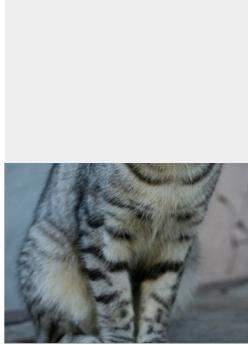


image  
completion



rotation  
prediction



"jigsaw puzzle"

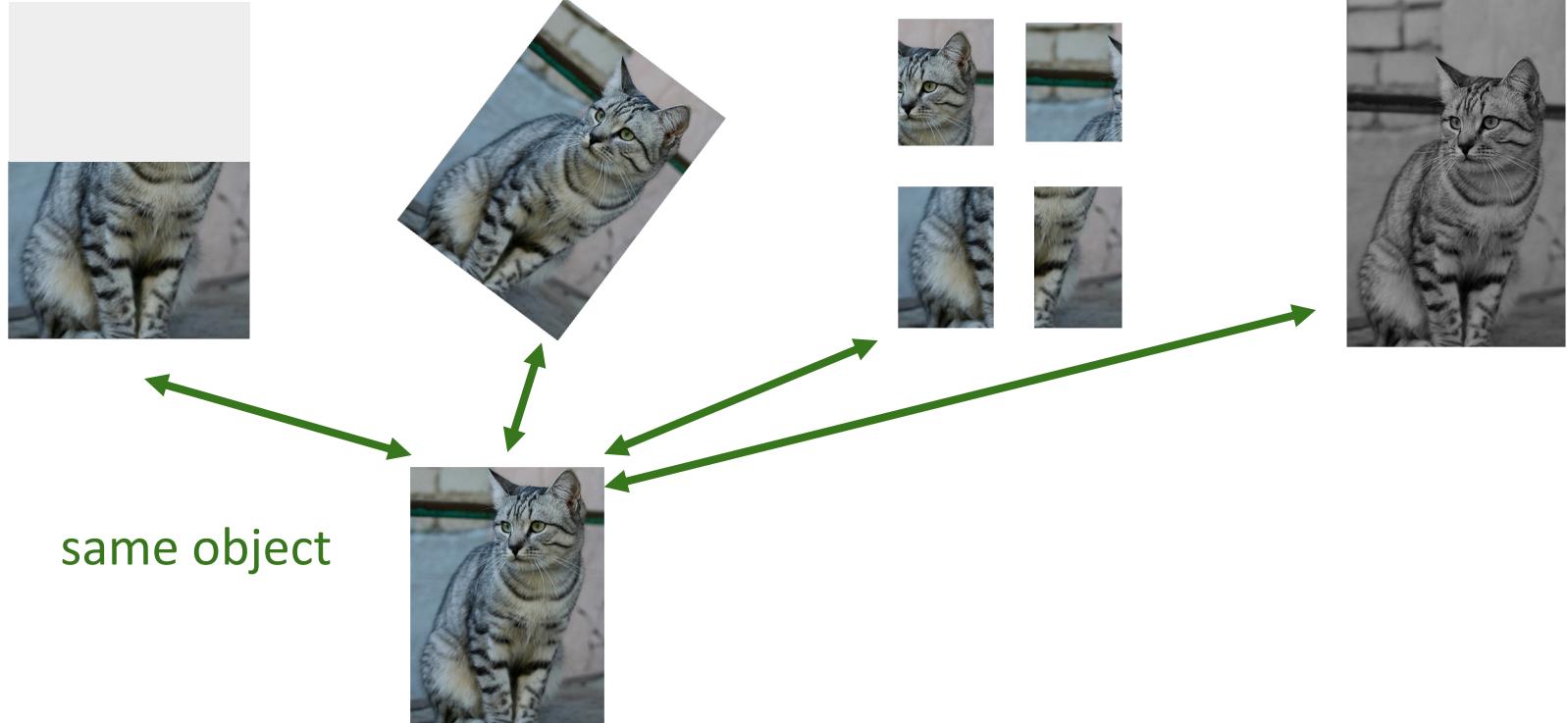


colorization

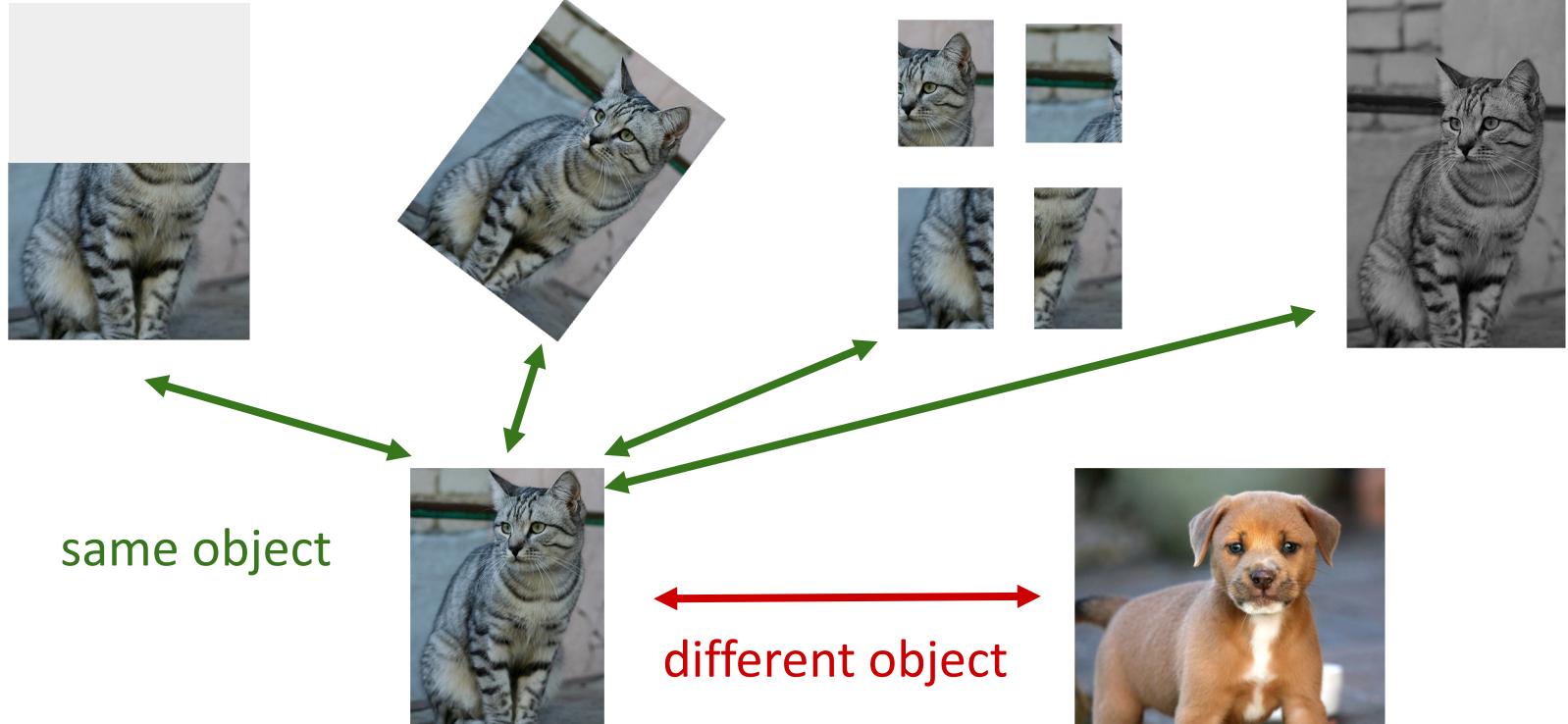
Learned representations may be tied to a specific pretext task!

Can we come up with a more general pretext task?

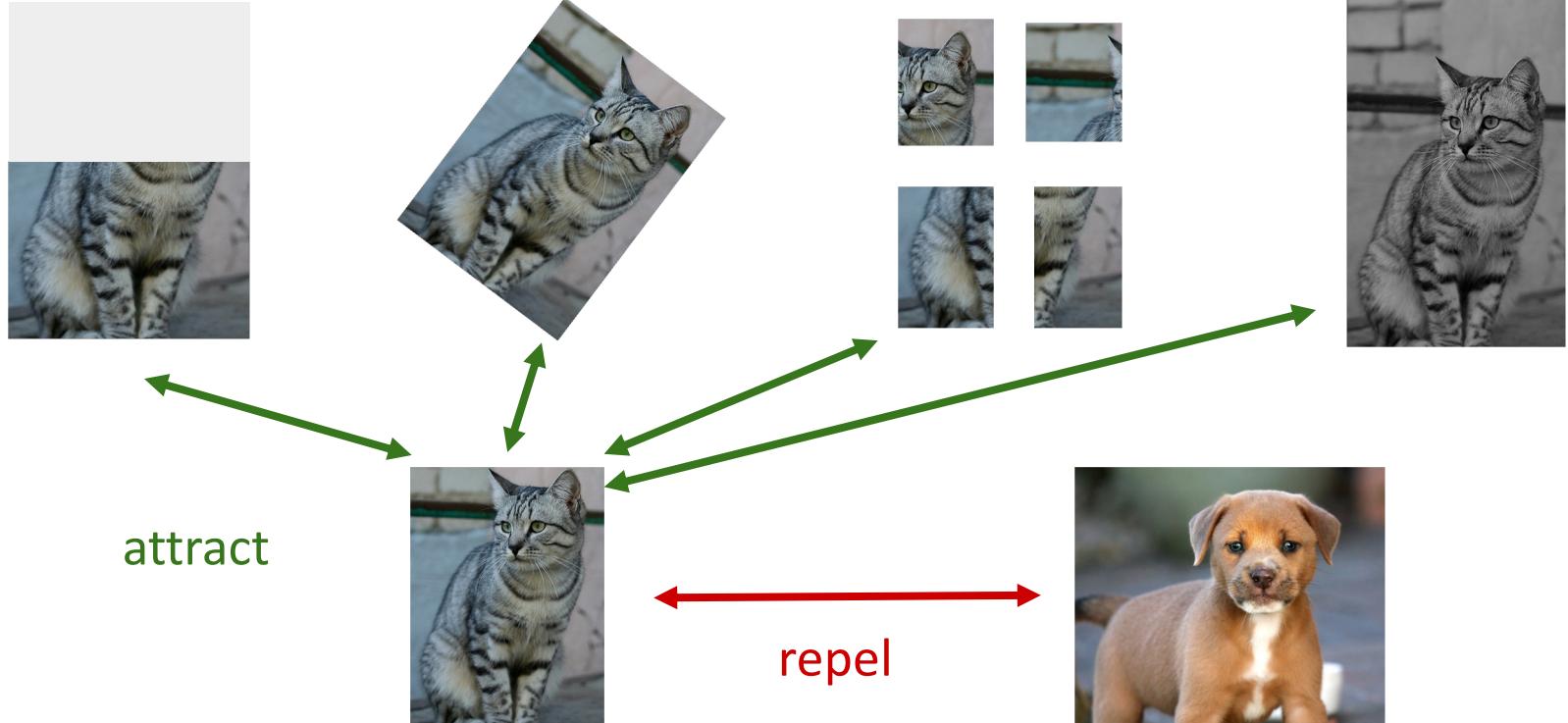
# A more general pretext task?



# A more general pretext task?



# Contrastive Representation Learning



# Today's Agenda

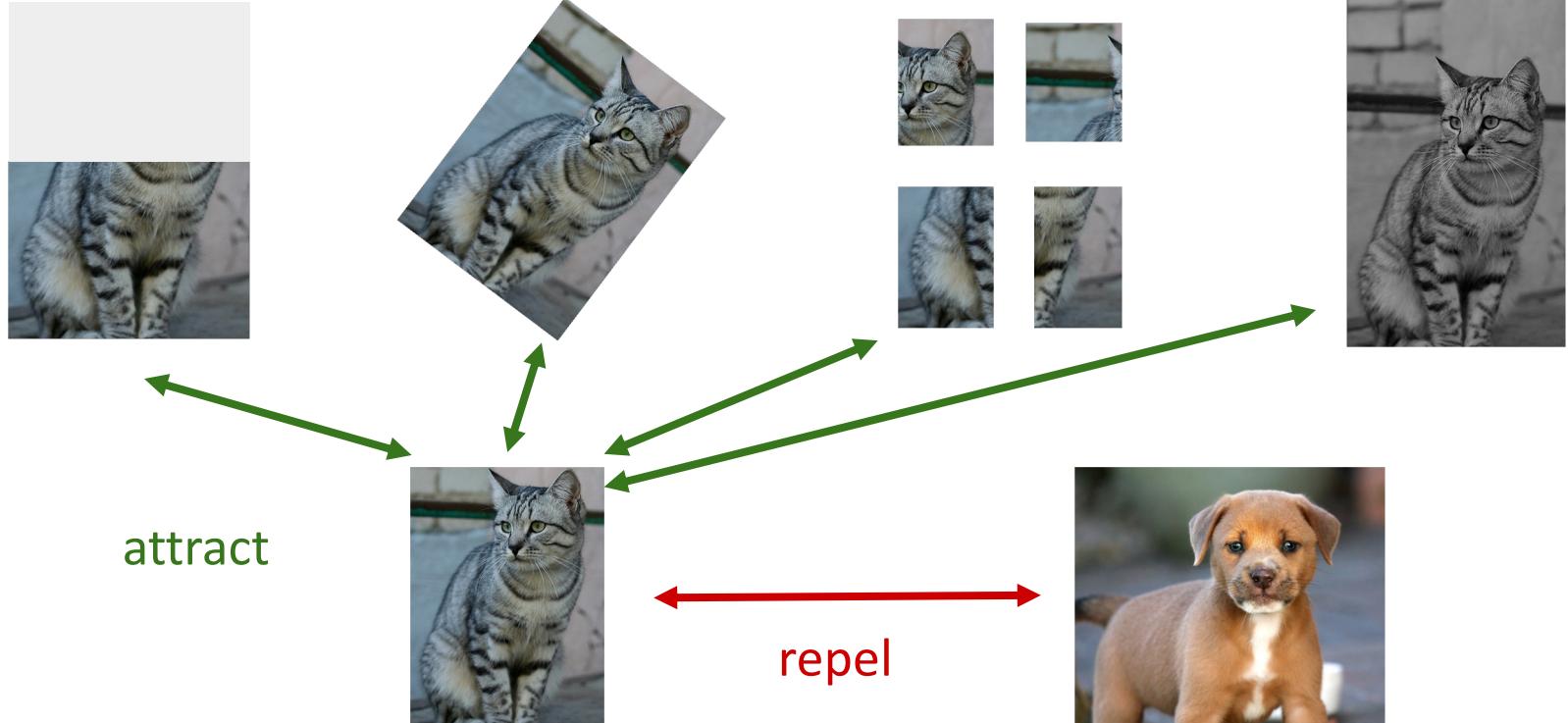
## Pretext tasks from image transformations

- Rotation, inpainting, rearrangement, coloring

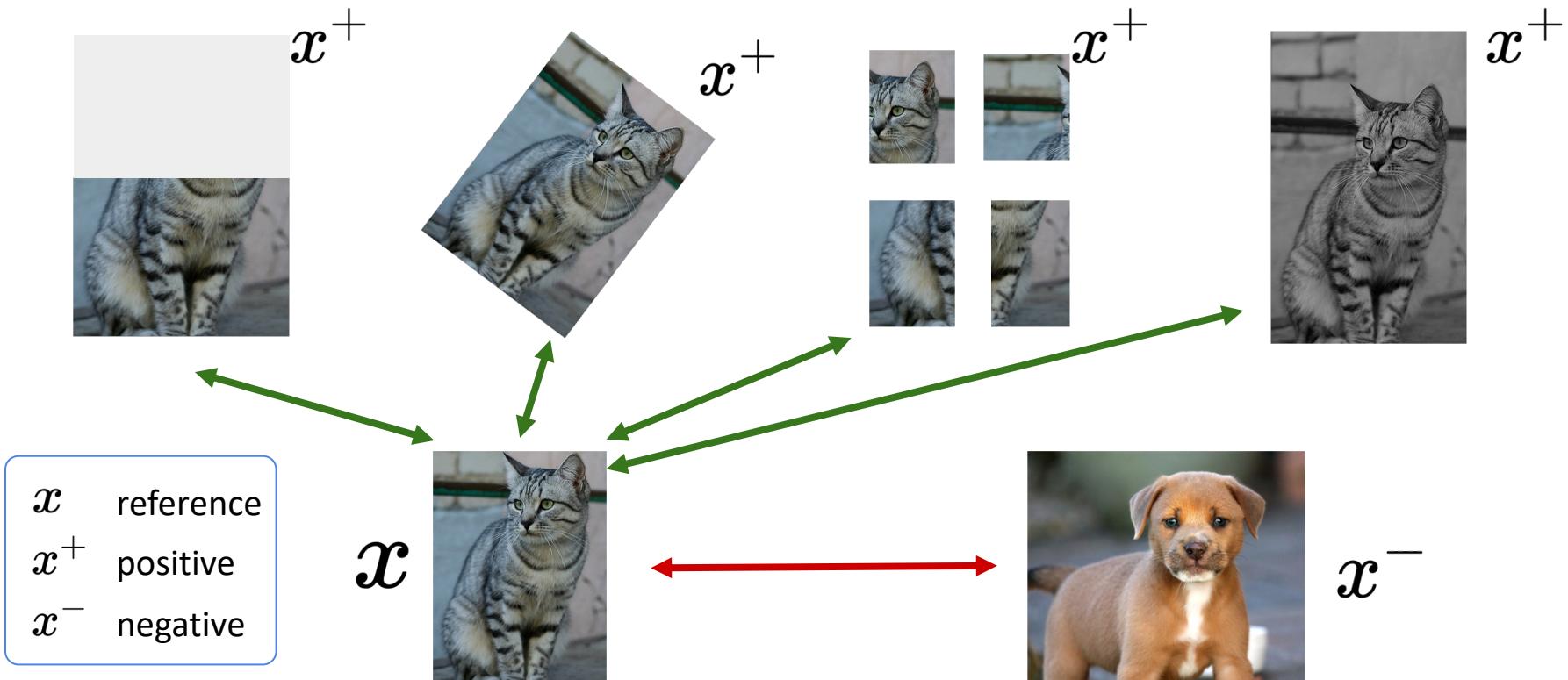
## Contrastive representation learning

- Intuition and formulation
- Instance contrastive learning: SimCLR and MOCO
- Sequence contrastive learning: CPC

# Contrastive Representation Learning



# Contrastive Representation Learning



# A formulation of contrastive learning

What we want:

$$\text{score}(f(x), f(x^+)) \gg \text{score}(f(x), f(x^-))$$

$x$ : reference sample;  $x^+$  positive sample;  $x^-$  negative sample

Given a chosen score function, we aim to learn an **encoder function**  $f$  that yields high score for positive pairs  $(x, x^+)$  and low scores for negative pairs  $(x, x^-)$ .

# A formulation of contrastive learning

Loss function given 1 positive sample and  $N - 1$  negative samples:

$$L = -\mathbb{E}_X \left[ \log \frac{\exp(s(f(x), f(x^+)))}{\exp(s(f(x), f(x^+))) + \sum_{j=1}^{N-1} \exp(s(f(x), f(x_j^-)))} \right]$$

# A formulation of contrastive learning

Loss function given 1 positive sample and  $N - 1$  negative samples:

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$x$



$x^+$



$x$



$x_1^-$



$x_2^-$



$x_3^-$

...

# A formulation of contrastive learning

Loss function given 1 positive sample and  $N - 1$  negative samples:

$$L = -\mathbb{E}_X \left[ \log \frac{\overline{\exp(s(f(x), f(x^+))}}}{\overline{\exp(s(f(x), f(x^+)) + \sum_{j=1}^{N-1} \exp(s(f(x), f(x_j^-)))}}} \right]$$

score for the positive pair                                    score for the N-1 negative pairs

This seems familiar ...

# A formulation of contrastive learning

Loss function given 1 positive sample and  $N - 1$  negative samples:

$$L = -\mathbb{E}_X \left[ \log \frac{\overline{\exp(s(f(x), f(x^+))}}}{\overline{\exp(s(f(x), f(x^+)) + \sum_{j=1}^{N-1} \exp(s(f(x), f(x_j^-)))}}} \right]$$

score for the positive pair                                    score for the N-1 negative pairs

This seems familiar ...

Cross entropy loss for a  $N$ -way softmax classifier!

i.e., learn to find the positive sample from the  $N$  samples

# A formulation of contrastive learning

Loss function given 1 positive sample and  $N - 1$  negative samples:

$$L = -\mathbb{E}_X \left[ \log \frac{\exp(s(f(x), f(x^+)))}{\exp(s(f(x), f(x^+))) + \sum_{j=1}^{N-1} \exp(s(f(x), f(x_j^-)))} \right]$$

Commonly known as the InfoNCE loss ([van den Oord et al., 2018](#))

A *lower bound* on the mutual information between  $f(x)$  and  $f(x^+)$

$$MI[f(x), f(x^+)] - \log(N) \geq -L$$

The larger the negative sample size ( $N$ ), the tighter the bound

Detailed derivation: [Poole et al., 2019](#)

# SimCLR: A Simple Framework for Contrastive Learning

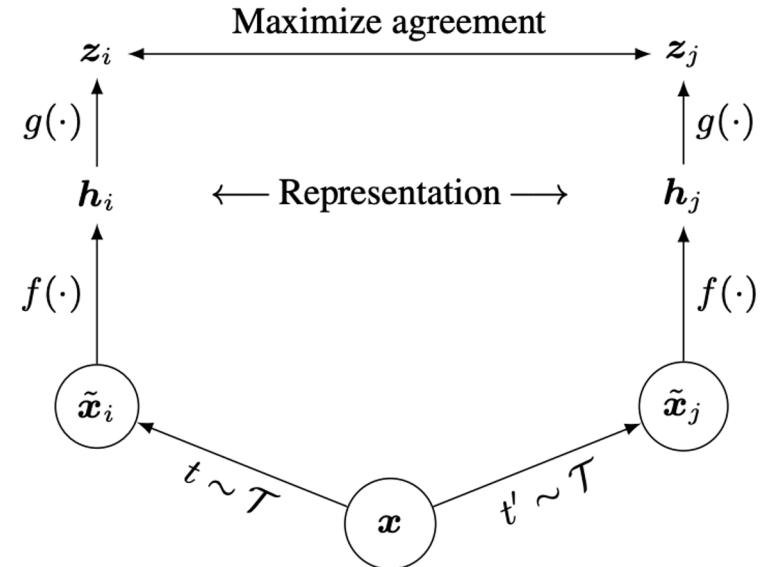
Cosine similarity as the score function:

$$s(u, v) = \frac{u^T v}{\|u\| \|v\|}$$

Use a projection network  $h(\cdot)$  to project features to a space where contrastive learning is applied

Generate positive samples through data augmentation:

- random cropping, random color distortion, and random blur.



Source: [Chen et al., 2020](#)

# SimCLR: generating positive samples from data augmentation



(a) Original



(b) Crop and resize



(c) Crop, resize (and flip)



(d) Color distort. (drop)



(e) Color distort. (jitter)



(f) Rotate  $\{90^\circ, 180^\circ, 270^\circ\}$



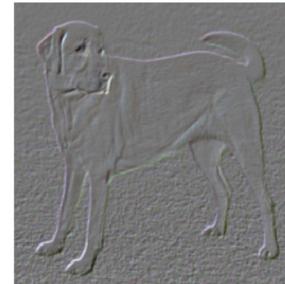
(g) Cutout



(h) Gaussian noise



(i) Gaussian blur



(j) Sobel filtering

Source: [Chen et al., 2020](#)

# SimCLR

Generate a positive pair  
by sampling data  
augmentation functions

---

**Algorithm 1** SimCLR's main learning algorithm.

---

```
input: batch size  $N$ , constant  $\tau$ , structure of  $f, g, \mathcal{T}$ .
for sampled minibatch  $\{\mathbf{x}_k\}_{k=1}^N$  do
    for all  $k \in \{1, \dots, N\}$  do
        draw two augmentation functions  $t \sim \mathcal{T}, t' \sim \mathcal{T}$ 
        # the first augmentation
         $\tilde{\mathbf{x}}_{2k-1} = t(\mathbf{x}_k)$ 
         $\mathbf{h}_{2k-1} = f(\tilde{\mathbf{x}}_{2k-1})$  # representation
         $\mathbf{z}_{2k-1} = g(\mathbf{h}_{2k-1})$  # projection
        # the second augmentation
         $\tilde{\mathbf{x}}_{2k} = t'(\mathbf{x}_k)$ 
         $\mathbf{h}_{2k} = f(\tilde{\mathbf{x}}_{2k})$  # representation
         $\mathbf{z}_{2k} = g(\mathbf{h}_{2k})$  # projection
    end for
    for all  $i \in \{1, \dots, 2N\}$  and  $j \in \{1, \dots, 2N\}$  do
         $s_{i,j} = \mathbf{z}_i^\top \mathbf{z}_j / (\|\mathbf{z}_i\| \|\mathbf{z}_j\|)$  # pairwise similarity
    end for
    define  $\ell(i, j)$  as  $\ell(i, j) = -\log \frac{\exp(s_{i,j}/\tau)}{\sum_{k=1}^{2N} \mathbb{1}_{[k \neq i]} \exp(s_{i,k}/\tau)}$ 
     $\mathcal{L} = \frac{1}{2N} \sum_{k=1}^N [\ell(2k-1, 2k) + \ell(2k, 2k-1)]$ 
    update networks  $f$  and  $g$  to minimize  $\mathcal{L}$ 
end for
return encoder network  $f(\cdot)$ , and throw away  $g(\cdot)$ 
```

---

Source: [Chen et al., 2020](#)

# SimCLR

Generate a positive pair  
by sampling data  
augmentation functions

---

**Algorithm 1** SimCLR's main learning algorithm.

---

```

input: batch size  $N$ , constant  $\tau$ , structure of  $f, g, \mathcal{T}$ .
for sampled minibatch  $\{\mathbf{x}_k\}_{k=1}^N$  do
    for all  $k \in \{1, \dots, N\}$  do
        draw two augmentation functions  $t \sim \mathcal{T}, t' \sim \mathcal{T}$ 
        # the first augmentation
         $\tilde{\mathbf{x}}_{2k-1} = t(\mathbf{x}_k)$ 
         $\mathbf{h}_{2k-1} = f(\tilde{\mathbf{x}}_{2k-1})$  # representation
         $\mathbf{z}_{2k-1} = g(\mathbf{h}_{2k-1})$  # projection
        # the second augmentation
         $\tilde{\mathbf{x}}_{2k} = t'(\mathbf{x}_k)$ 
         $\mathbf{h}_{2k} = f(\tilde{\mathbf{x}}_{2k})$  # representation
         $\mathbf{z}_{2k} = g(\mathbf{h}_{2k})$  # projection
    end for
    for all  $i \in \{1, \dots, 2N\}$  and  $j \in \{1, \dots, 2N\}$  do
         $s_{i,j} = \mathbf{z}_i^\top \mathbf{z}_j / (\|\mathbf{z}_i\| \|\mathbf{z}_j\|)$  # pairwise similarity
    end for
    define  $\ell(i, j)$  as  $\ell(i, j) = -\log \frac{\exp(s_{i,j}/\tau)}{\sum_{k=1}^{2N} \mathbb{1}_{[k \neq i]} \exp(s_{i,k}/\tau)}$ 
     $\mathcal{L} = \frac{1}{2N} \sum_{k=1}^N [\ell(2k-1, 2k) + \ell(2k, 2k-1)]$ 
    update networks  $f$  and  $g$  to minimize  $\mathcal{L}$ 
end for
return encoder network  $f(\cdot)$ , and throw away  $g(\cdot)$ 

```

---

InfoNCE loss:  
Use all non-positive  
samples in the batch  
as  $\mathbf{x}^-$

Source: [Chen et al.,  
2020](#)

# SimCLR

Generate a positive pair  
by sampling data  
augmentation functions

Iterate through and use  
each of the  $2N$  sample as  
reference, compute  
average loss

---

**Algorithm 1** SimCLR's main learning algorithm.

---

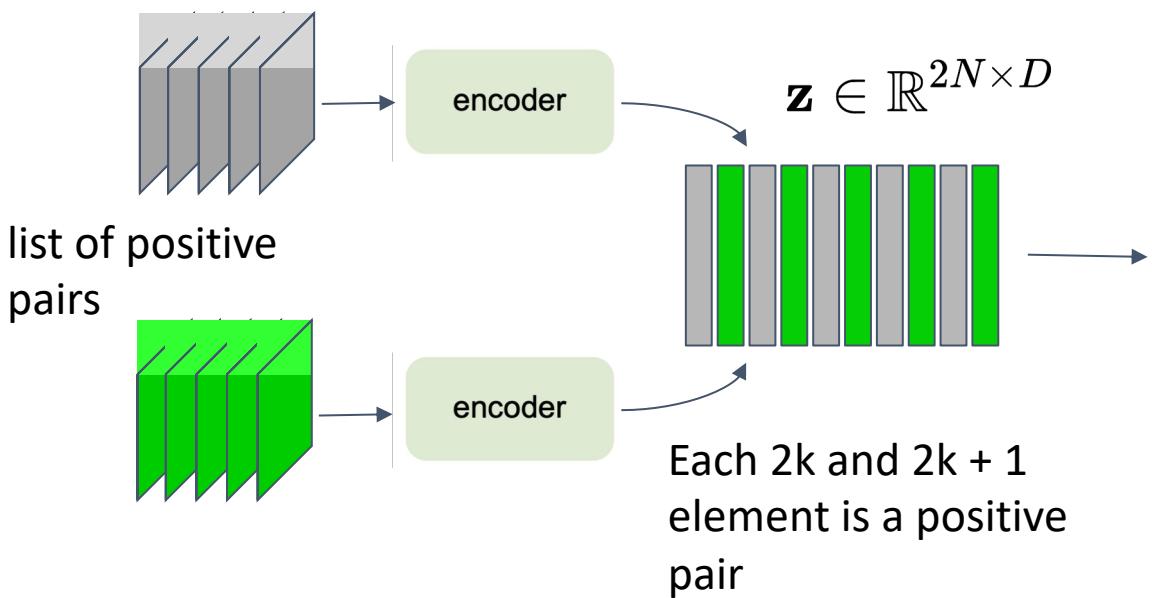
```
input: batch size  $N$ , constant  $\tau$ , structure of  $f, g, \mathcal{T}$ .
for sampled minibatch  $\{\mathbf{x}_k\}_{k=1}^N$  do
    for all  $k \in \{1, \dots, N\}$  do
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         $\mathbf{z}_{2k-1} = g(\mathbf{h}_{2k-1})$  # projection
        # the second augmentation
         $\tilde{\mathbf{x}}_{2k} = t'(\mathbf{x}_k)$ 
         $\mathbf{h}_{2k} = f(\tilde{\mathbf{x}}_{2k})$  # representation
         $\mathbf{z}_{2k} = g(\mathbf{h}_{2k})$  # projection
    end for
    for all  $i \in \{1, \dots, 2N\}$  and  $j \in \{1, \dots, 2N\}$  do
         $s_{i,j} = \mathbf{z}_i^\top \mathbf{z}_j / (\|\mathbf{z}_i\| \|\mathbf{z}_j\|)$  # pairwise similarity
    end for
    define  $\ell(i, j)$  as  $\ell(i, j) = -\log \frac{\exp(s_{i,j}/\tau)}{\sum_{k=1}^{2N} \mathbb{1}_{[k \neq i]} \exp(s_{i,k}/\tau)}$ 
     $\mathcal{L} = \frac{1}{2N} \sum_{k=1}^N [\ell(2k-1, 2k) + \ell(2k, 2k-1)]$ 
    update networks  $f$  and  $g$  to minimize  $\mathcal{L}$ 
end for
return encoder network  $f(\cdot)$ , and throw away  $g(\cdot)$ 
```

---

InfoNCE loss:  
Use all non-positive  
samples in the batch  
as  $\mathbf{x}^-$

Source: [Chen et al.,  
2020](#)

# SimCLR: mini-batch training



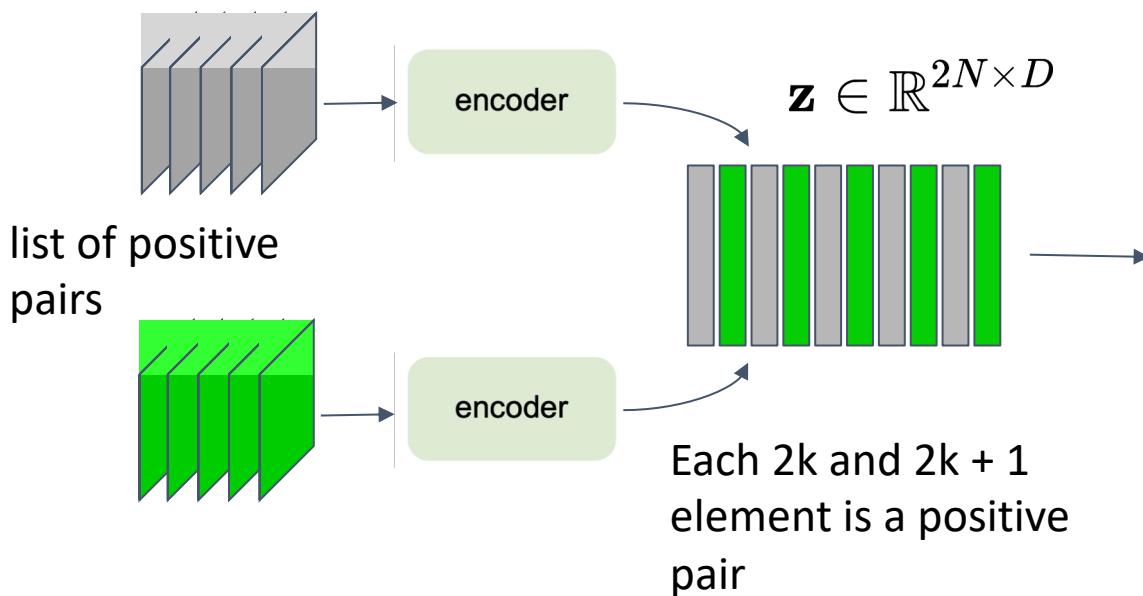
$$s_{i,j} = \frac{\mathbf{z}_i^T \mathbf{z}_j}{\|\mathbf{z}_i\| \|\mathbf{z}_j\|}$$

"Affinity matrix"


$2N$

$2N$

# SimCLR: mini-batch training



$$s_{i,j} = \frac{\mathbf{z}_i^T \mathbf{z}_j}{\|\mathbf{z}_i\| \|\mathbf{z}_j\|}$$

"Affinity matrix"

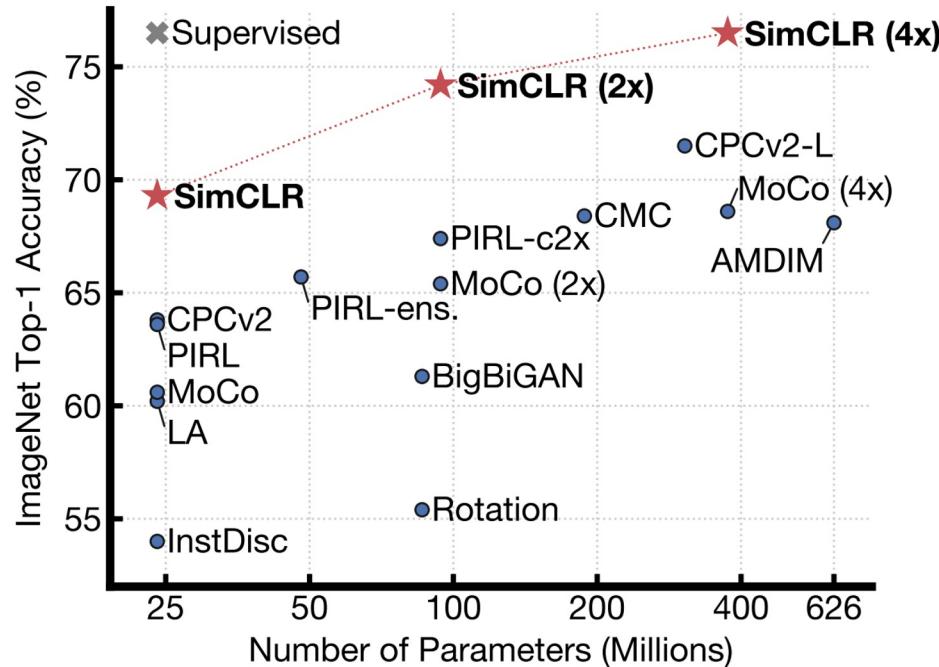
A 2D grid representing the "Affinity matrix". The grid has  $2N$  rows and  $2N$  columns, indicated by the labels "2N" at the bottom right. The grid contains several blue square blocks of size 2x2. These blue blocks are located at positions (0,1), (1,0), (2,3), (3,2), (4,5), (5,4), (6,7), and (7,6). All other cells in the grid are white.

$2N$



= classification label for each row

# Training linear classifier on SimCLR features



Train feature encoder on **ImageNet** (entire training set) using SimCLR.

Freeze feature encoder, train a linear classifier on top with labeled data.

Source: [Chen et al., 2020](#)

# Semi-supervised learning on SimCLR features

Method	Architecture	Label fraction		
		1%	10%	Top 5
Supervised baseline	ResNet-50	48.4	80.4	
<i>Methods using other label-propagation:</i>				
Pseudo-label	ResNet-50	51.6	82.4	
VAT+Entropy Min.	ResNet-50	47.0	83.4	
UDA (w. RandAug)	ResNet-50	-	88.5	
FixMatch (w. RandAug)	ResNet-50	-	89.1	
S4L (Rot+VAT+En. M.)	ResNet-50 (4×)	-	91.2	
<i>Methods using representation learning only:</i>				
InstDisc	ResNet-50	39.2	77.4	
BigBiGAN	RevNet-50 (4×)	55.2	78.8	
PIRL	ResNet-50	57.2	83.8	
CPC v2	ResNet-161(*)	77.9	91.2	
SimCLR (ours)	ResNet-50	75.5	87.8	
SimCLR (ours)	ResNet-50 (2×)	83.0	91.2	
SimCLR (ours)	ResNet-50 (4×)	<b>85.8</b>	<b>92.6</b>	

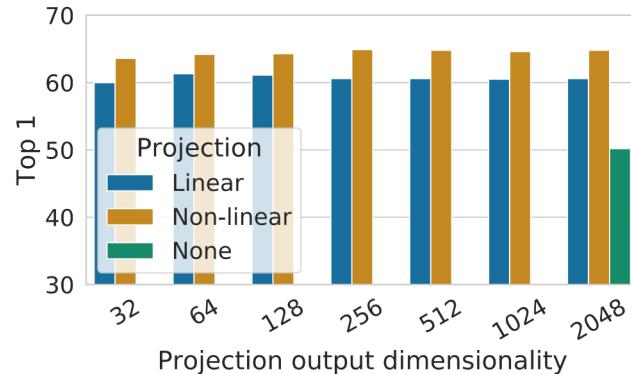
Table 7. ImageNet accuracy of models trained with few labels.

Train feature encoder on **ImageNet** (entire training set) using SimCLR.

**Finetune** the encoder with 1% / 10% of labeled data on ImageNet.

Source: [Chen et al., 2020](#)

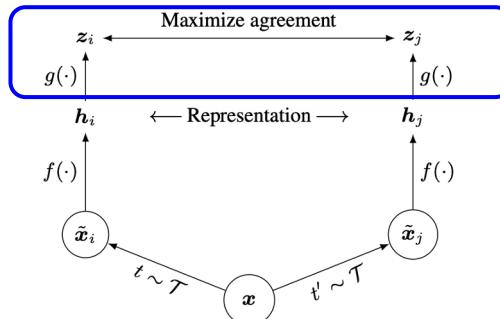
# SimCLR design choices: projection head



Linear / non-linear projection heads improve representation learning.

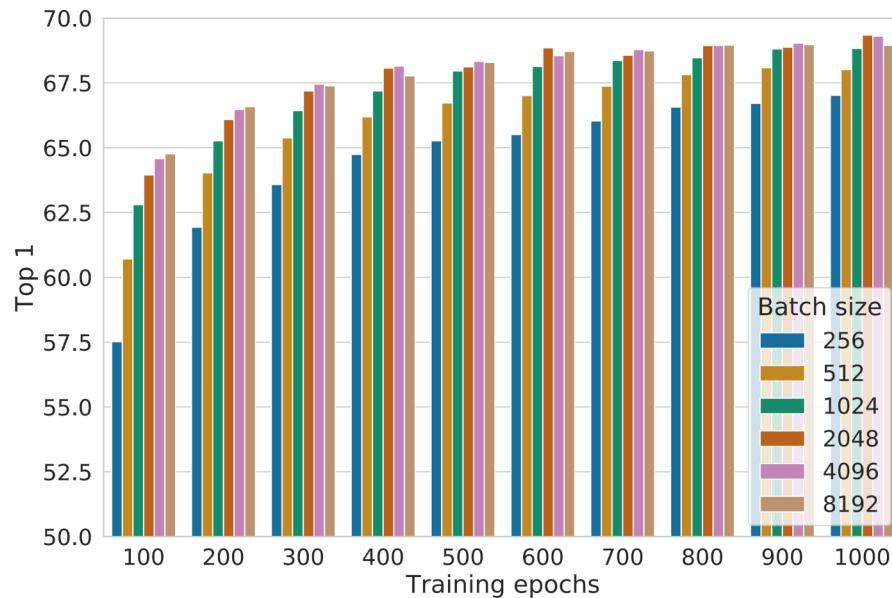
A possible explanation:

- contrastive learning objective may discard useful information for downstream tasks
- representation space  $\mathbf{z}$  is trained to be invariant to data transformation.
- by leveraging the projection head  $\mathbf{g}(\cdot)$ , more information can be preserved in the  $\mathbf{h}$  representation space



Source: [Chen et al., 2020](#)

# SimCLR design choices: large batch size



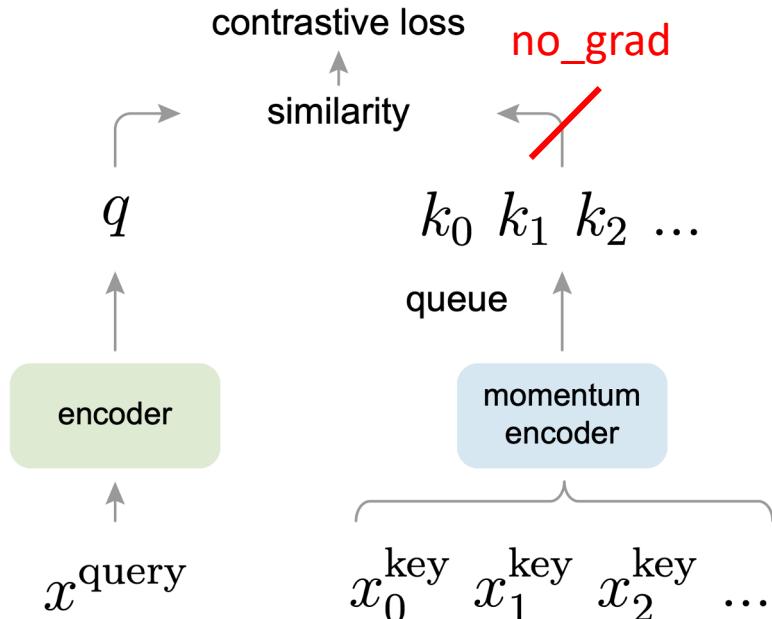
Large training batch size is crucial for SimCLR!

Large batch size causes large memory footprint during backpropagation:  
requires distributed training on TPUs  
(ImageNet experiments)

Figure 9. Linear evaluation models (ResNet-50) trained with different batch size and epochs. Each bar is a single run from scratch.<sup>10</sup>

Source: [Chen et al., 2020](#)

# Momentum Contrastive Learning (MoCo)

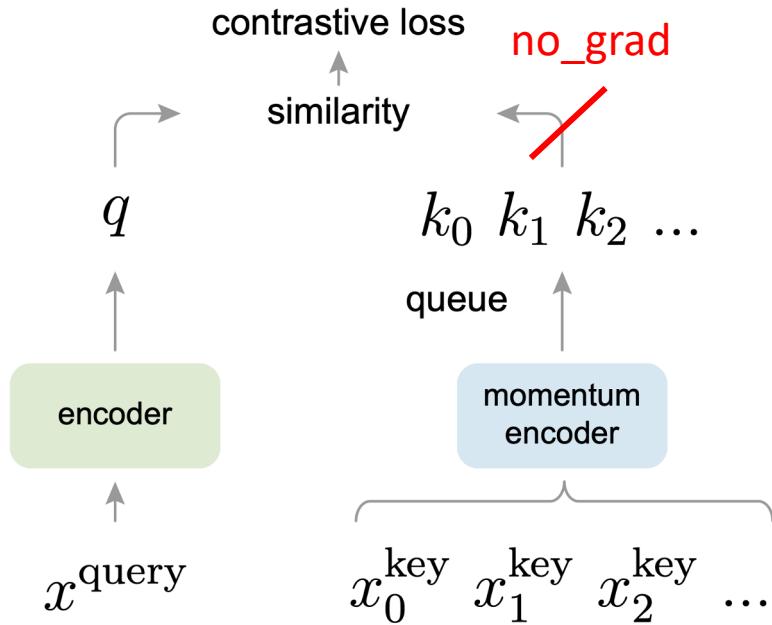


## Key differences to SimCLR:

- Keep a running **queue** of keys (negative samples).
- Compute gradients and update the encoder **only through the queries**.
- Decouple min-batch size with the number of keys: can support **a large number of negative samples**.

Source: [He et al., 2020](#)

# Momentum Contrastive Learning (MoCo)



## Key differences to SimCLR:

- Keep a running **queue** of keys (negative samples).
- Compute gradients and update the encoder **only through the queries**.
- Decouple min-batch size with the number of keys: can support **a large number of negative samples**.
- The key encoder is **slowly progressing** through the momentum update rules:  
$$\theta_k \leftarrow m\theta_k + (1 - m)\theta_q$$

Source: [He et al., 2020](#)

# MoCo

Generate a positive pair  
by sampling data  
augmentation functions

No gradient through  
the positive sample

Update the FIFO negative  
sample queue

**Algorithm 1** Pseudocode of MoCo in a PyTorch-like style.

```
# f_q, f_k: encoder networks for query and key
# queue: dictionary as a queue of K keys (CxK)
# m: momentum
# t: temperature

f_k.params = f_q.params # initialize
for x in loader: # load a minibatch x with N samples
    x_q = aug(x) # a randomly augmented version
    x_k = aug(x) # another randomly augmented version

    q = f_q.forward(x_q) # queries: NxC
    k = f_k.forward(x_k) # keys: NxC
    k = k.detach() # no gradient to keys

    # positive logits: Nx1
    l_pos = bmm(q.view(N, 1, C), k.view(N, C, 1))

    # negative logits: NxK
    l_neg = mm(q.view(N, C), queue.view(C, K))

    # logits: Nx(1+K)
    logits = cat([l_pos, l_neg], dim=1)

    # contrastive loss, Eqn.(1)
    labels = zeros(N) # positives are the 0-th
    loss = CrossEntropyLoss(logits/t, labels)

    # SGD update: query network
    loss.backward()
    update(f_q.params)

    # momentum update: key network
    f_k.params = m*f_k.params+(1-m)*f_q.params

    # update dictionary
    enqueue(queue, k) # enqueue the current minibatch
    dequeue(queue) # dequeue the earliest minibatch
```

bmm: batch matrix multiplication; mm: matrix multiplication; cat: concatenation.

Use the running queue  
of keys as the negative  
samples

InfoNCE loss

Update  $f_k$  through  
momentum

Source: [He et al., 2020](#)

# “MoCo V2”

## Improved Baselines with Momentum Contrastive Learning

Xinlei Chen Haoqi Fan Ross Girshick Kaiming He  
Facebook AI Research (FAIR)

A hybrid of ideas from SimCLR and MoCo:

- **From SimCLR:** non-linear projection head and strong data augmentation.
- **From MoCo:** momentum-updated queues that allow training on a large number of negative samples (no TPU required!).

Source: [Chen et al., 2020](#)

# MoCo vs. SimCLR vs. MoCo V2

case	unsup. pre-train				ImageNet acc.	VOC detection		
	MLP	aug+	cos	epochs		AP <sub>50</sub>	AP	AP <sub>75</sub>
supervised					76.5	81.3	53.5	58.8
MoCo v1				200	60.6	81.5	55.9	62.6
(a)	✓			200	66.2	82.0	56.4	62.6
(b)		✓		200	63.4	82.2	56.8	63.2
(c)	✓	✓		200	67.3	<b>82.5</b>	57.2	63.9
(d)	✓	✓	✓	200	67.5	82.4	57.0	63.6
(e)	✓	✓	✓	<b>800</b>	<b>71.1</b>	<b>82.5</b>	<b>57.4</b>	<b>64.0</b>

Table 1. **Ablation of MoCo baselines**, evaluated by ResNet-50 for (i) ImageNet linear classification, and (ii) fine-tuning VOC object detection (mean of 5 trials). “MLP”: with an MLP head; “aug+”: with extra blur augmentation; “cos”: cosine learning rate schedule.

## Key takeaways:

- Non-linear projection head and strong data augmentation are crucial for contrastive learning.

Source: [Chen et al., 2020](#)

# MoCo vs. SimCLR vs. MoCo V2

case	MLP	aug+	cos	unsup. pre-train epochs	batch	ImageNet acc.
MoCo v1 [6]				200	256	60.6
SimCLR [2]	✓	✓	✓	200	256	61.9
SimCLR [2]	✓	✓	✓	200	8192	66.6
<b>MoCo v2</b>	✓	✓	✓	200	256	<b>67.5</b>

*results of longer unsupervised training follow:*

SimCLR [2]	✓	✓	✓	1000	4096	69.3
<b>MoCo v2</b>	✓	✓	✓	800	256	<b>71.1</b>

Table 2. **MoCo vs. SimCLR**: ImageNet linear classifier accuracy (**ResNet-50, 1-crop 224×224**), trained on features from unsupervised pre-training. “aug+” in SimCLR includes blur and stronger color distortion. SimCLR ablations are from Fig. 9 in [2] (we thank the authors for providing the numerical results).

## Key takeaways:

- Non-linear projection head and strong data augmentation are crucial for contrastive learning.
- Decoupling mini-batch size with negative sample size allows MoCo-V2 to outperform SimCLR with smaller batch size (256 vs. 8192).

Source: [Chen et al., 2020](#)

# MoCo vs. SimCLR vs. MoCo V2

mechanism	batch	memory / GPU	time / 200-ep.
MoCo	256	<b>5.0G</b>	<b>53 hrs</b>
end-to-end	256	7.4G	65 hrs
end-to-end	4096	93.0G <sup>†</sup>	n/a

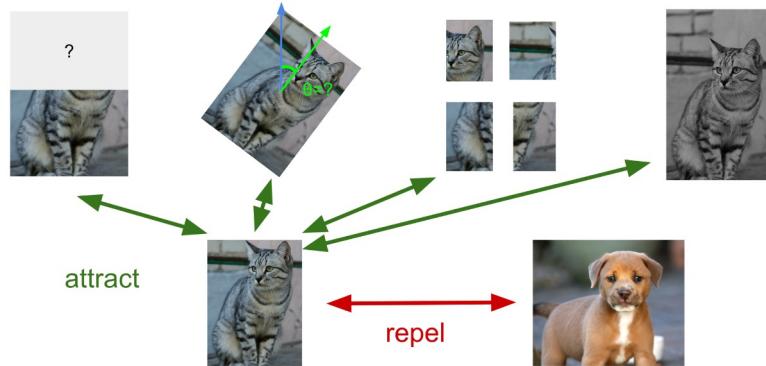
Table 3. **Memory and time cost** in 8 V100 16G GPUs, implemented in PyTorch. <sup>†</sup>: based on our estimation.

## Key takeaways:

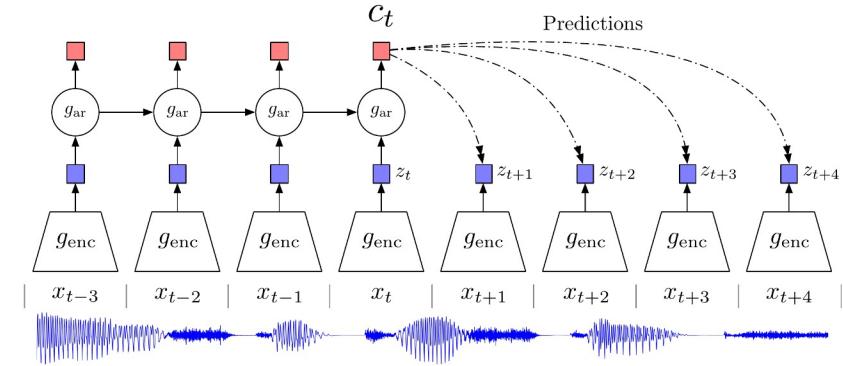
- Non-linear projection head and strong data augmentation are crucial for contrastive learning.
- Decoupling mini-batch size with negative sample size allows MoCo-V2 to outperform SimCLR with smaller batch size (256 vs. 8192).
- ... all with much smaller memory footprint! (“end-to-end” means SimCLR here)

Source: [Chen et al., 2020](#)

# Instance vs. Sequence Contrastive Learning



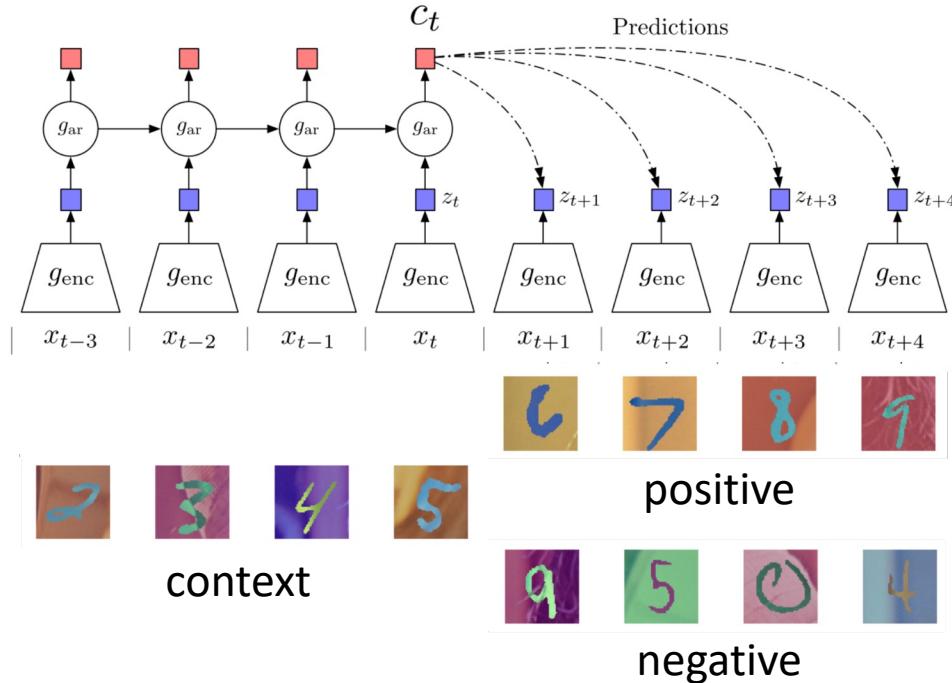
**Instance-level contrastive learning:**  
contrastive learning based on  
**positive & negative instances.**  
Examples: SimCLR, MoCo



Source: [van den Oord et al., 2018](#)

**Sequence-level contrastive learning:**  
contrastive learning based on  
**sequential / temporal orders.**  
Example: **Contrastive Predictive Coding (CPC)**

# Contrastive Predictive Coding (CPC)



**Contrastive:** contrast between “right” and “wrong” sequences using contrastive learning.

**Predictive:** the model has to predict future patterns given the current context.

**Coding:** the model learns useful feature vectors, or “code”, for downstream tasks, similar to other self-supervised methods.

Figure [source](#)

Source: [van den Oord et al., 2018](#),

# Contrastive Predictive Coding (CPC)

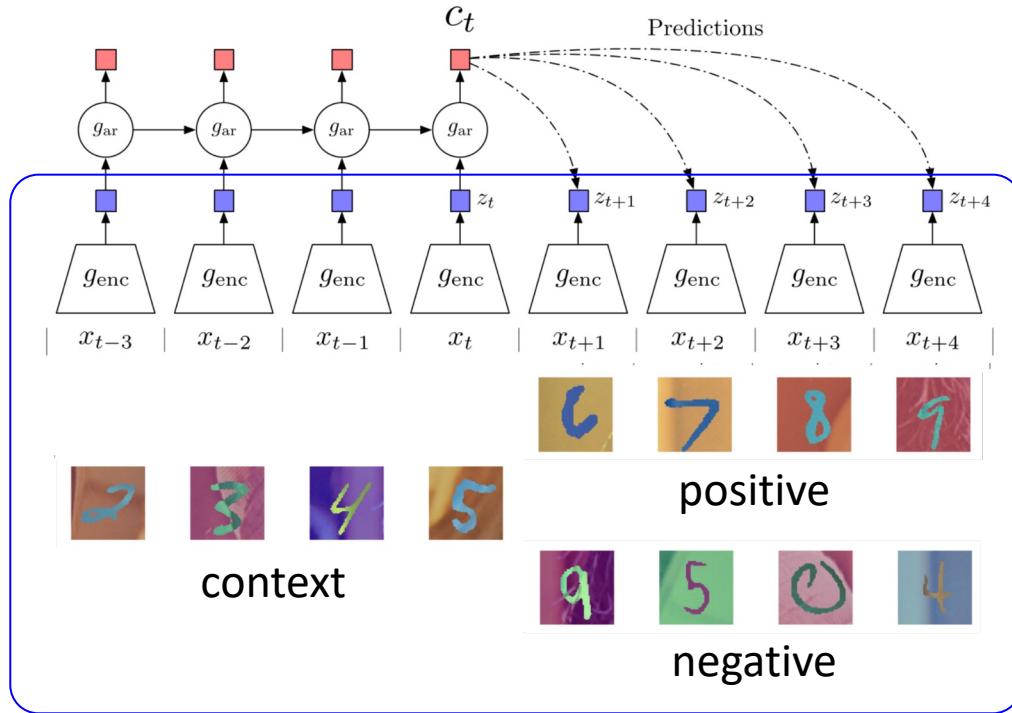


Figure [source](#)

1. Encode all samples in a sequence into vectors  $z_t = g_{enc}(x_t)$

Source: [van den Oord et al., 2018](#),

# Contrastive Predictive Coding (CPC)

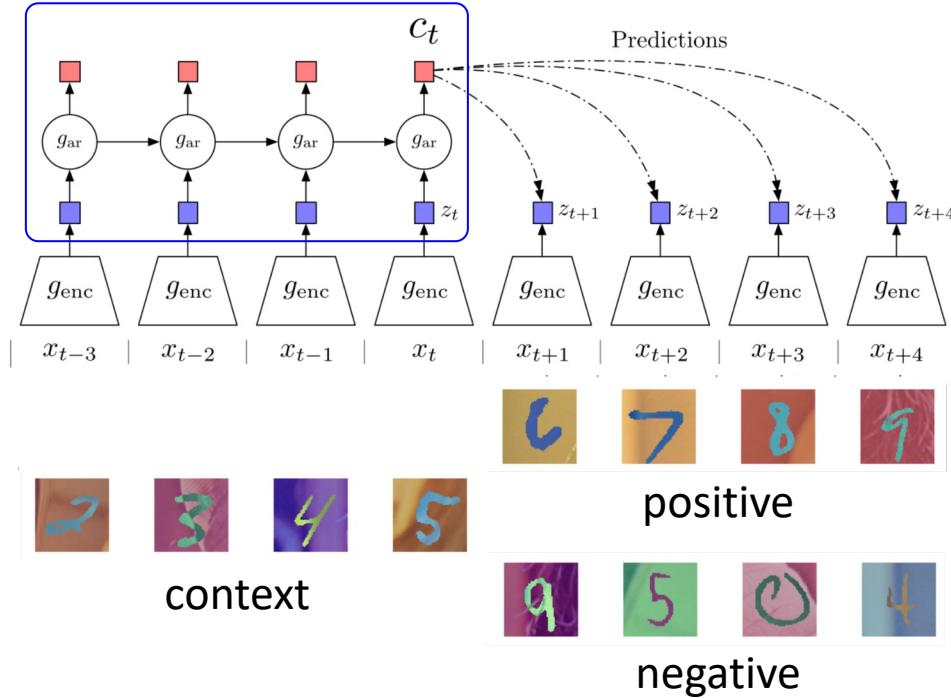


Figure [source](#)

1. Encode all samples in a sequence into vectors  $\mathbf{z}_t = \mathbf{g}_{enc}(\mathbf{x}_t)$
2. Summarize context (e.g., half of a sequence) into a context code  $\mathbf{c}_t$  using an auto-regressive model ( $\mathbf{g}_{ar}$ ).

Source: [van den Oord et al., 2018](#),

# Contrastive Predictive Coding (CPC)

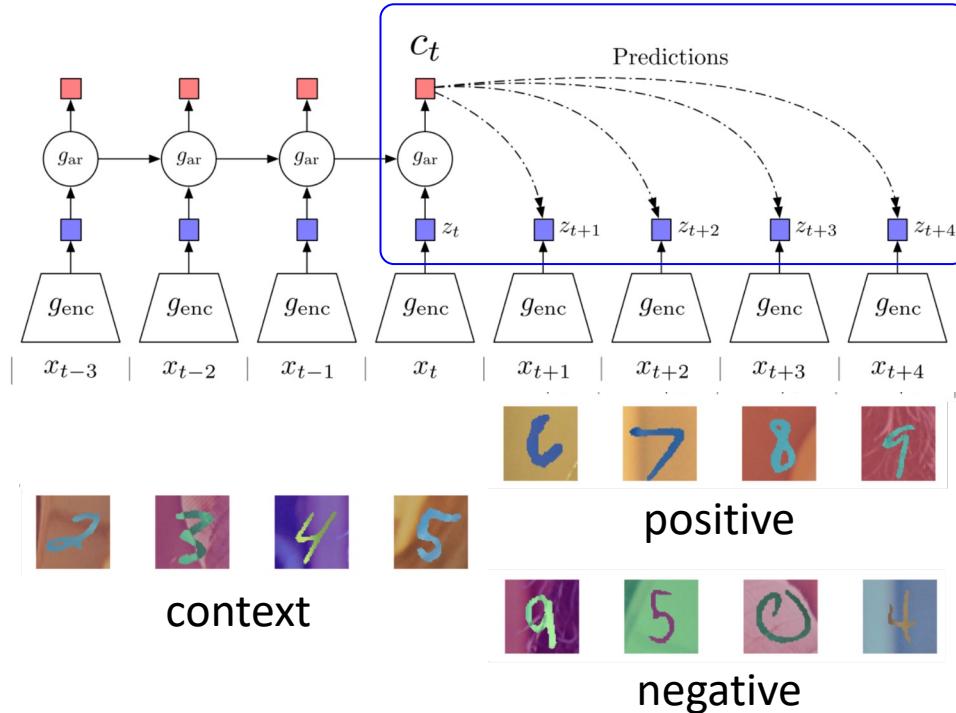


Figure [source](#)

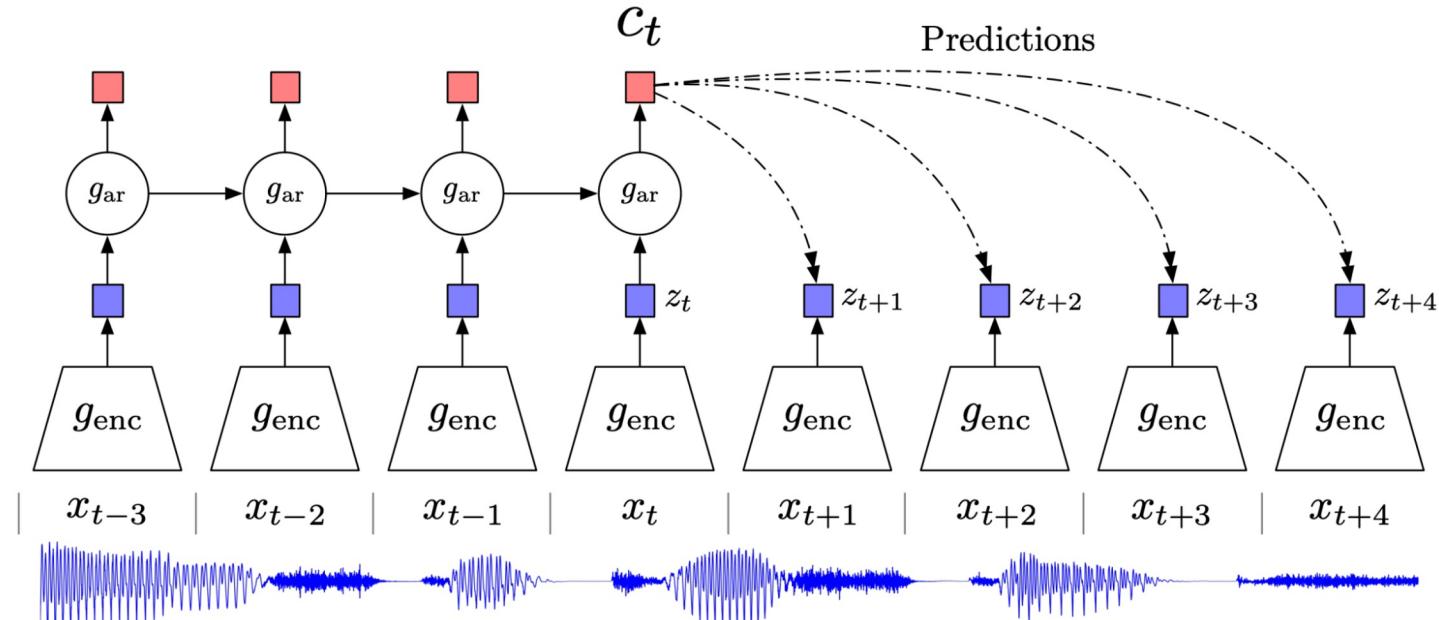
1. Encode all samples in a sequence into vectors  $\mathbf{z}_t = \mathbf{g}_{enc}(\mathbf{x}_t)$
2. Summarize context (e.g., half of a sequence) into a context code  $\mathbf{c}_t$  using an auto-regressive model ( $\mathbf{g}_{ar}$ ).
3. Compute InfoNCE loss between the context  $\mathbf{c}_t$  and future code  $\mathbf{z}_{t+k}$  using the following **time-dependent score function**:

$$s_k(z_{t+k}, c_t) = z_{t+k}^T W_k c_t$$

, where  $W_k$  is a trainable matrix.

Source: [van den Oord et al., 2018](#),

# CPC example: modeling audio sequences



Source: [van den Oord et al., 2018](#),

# CPC example: modeling audio sequences

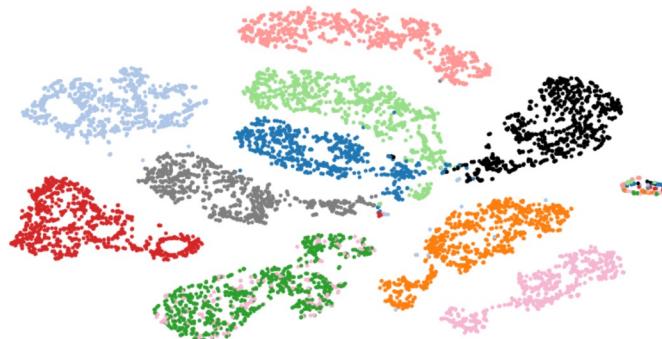


Figure 2: t-SNE visualization of audio (speech) representations for a subset of 10 speakers (out of 251). Every color represents a different speaker.

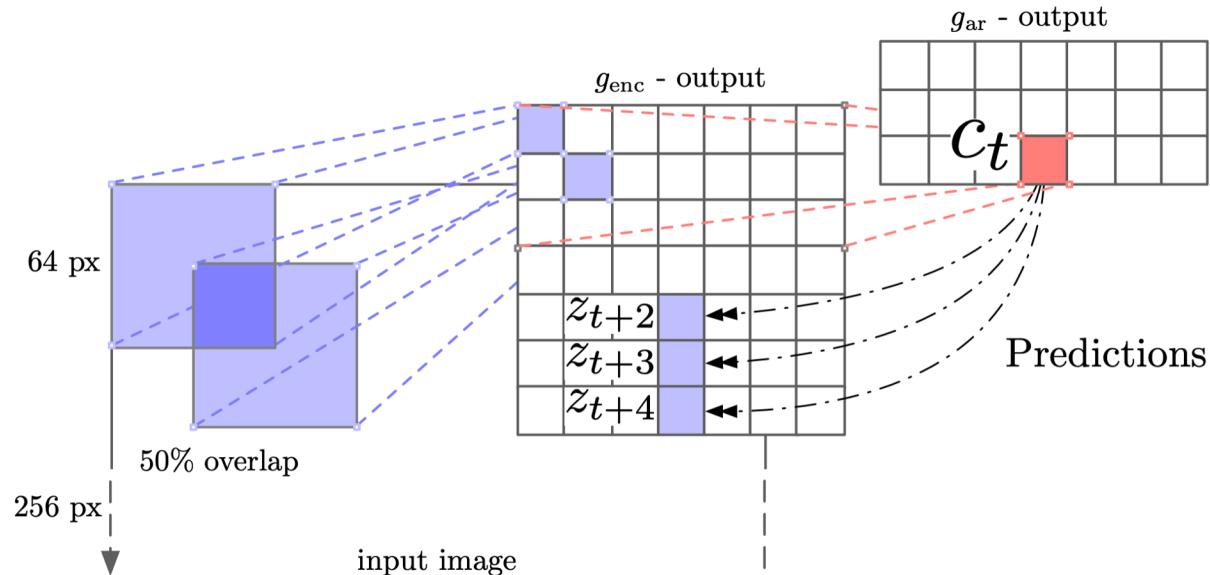
Method	ACC
<b>Phone classification</b>	
Random initialization	27.6
MFCC features	39.7
CPC	64.6
Supervised	74.6
<b>Speaker classification</b>	
Random initialization	1.87
MFCC features	17.6
CPC	97.4
Supervised	98.5

Linear classification on trained representations (LibriSpeech dataset)

Source: [van den Oord et al., 2018](#),

# CPC example: modeling visual context

**Idea:** split image into patches, model rows of patches from top to bottom as a sequence. I.e., use top rows as context to predict bottom rows.



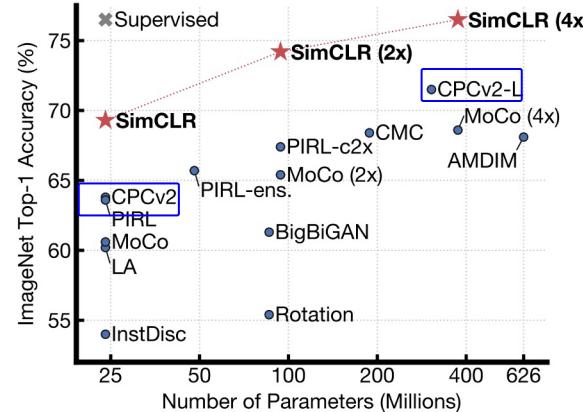
Source: [van den Oord et al., 2018](#),

# CPC example: modeling visual context

Method	Top-1 ACC
<b>Using AlexNet conv5</b>	
Video [28]	29.8
Relative Position [11]	30.4
BiGan [35]	34.8
Colorization [10]	35.2
Jigsaw [29] *	38.1
<b>Using ResNet-V2</b>	
Motion Segmentation [36]	27.6
Exemplar [36]	31.5
Relative Position [36]	36.2
Colorization [36]	39.6
<b>CPC</b>	<b>48.7</b>

Table 3: ImageNet top-1 unsupervised classification results. \*Jigsaw is not directly comparable to the other AlexNet results because of architectural differences.

- Compares favorably with other pretext task-based self-supervised learning method.
- Doesn't do as well compared to newer instance-based contrastive learning methods on image feature learning.



Source: [van den Oord et al., 2018](#),

# Summary: Contrastive Representation Learning

A general formulation for contrastive learning:

$$\text{score}(f(x), f(x^+)) \gg \text{score}(f(x), f(x^-))$$

InfoNCE loss: N-way classification among positive and negative samples

$$L = -\mathbb{E}_X \left[ \log \frac{\exp(s(f(x), f(x^+)))}{\exp(s(f(x), f(x^+))) + \sum_{j=1}^{N-1} \exp(s(f(x), f(x_j^-)))} \right]$$

Commonly known as the InfoNCE loss ([van den Oord et al., 2018](#))

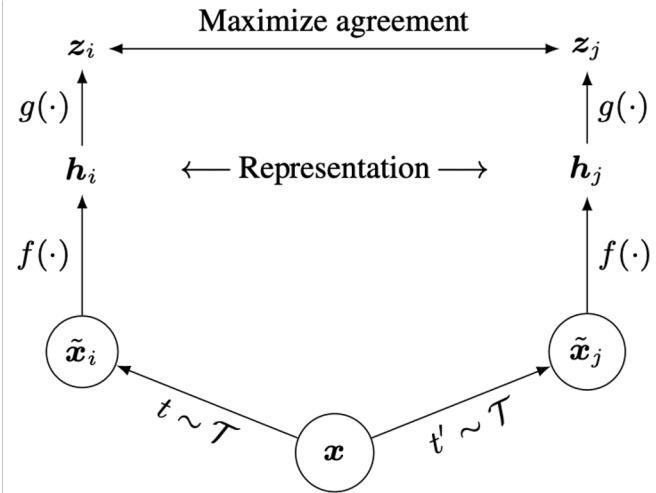
A *lower bound* on the mutual information between  $f(x)$  and  $f(x^+)$

$$MI[f(x), f(x^+)] - \log(N) \geq -L$$

# Summary: Contrastive Representation Learning

**SimCLR**: a simple framework for contrastive representation learning

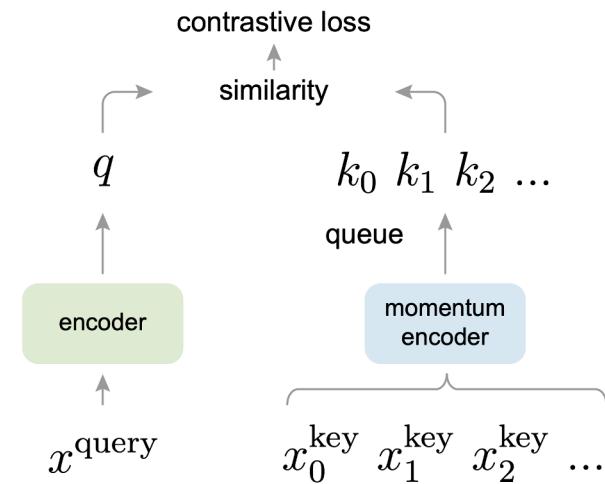
- **Key ideas**: non-linear projection head to allow flexible representation learning
- Simple to implement, effective in learning visual representation
- Requires large training batch size to be effective; large memory footprint



# Summary: Contrastive Representation Learning

**MoCo (v1, v2):** contrastive learning using momentum sample encoder

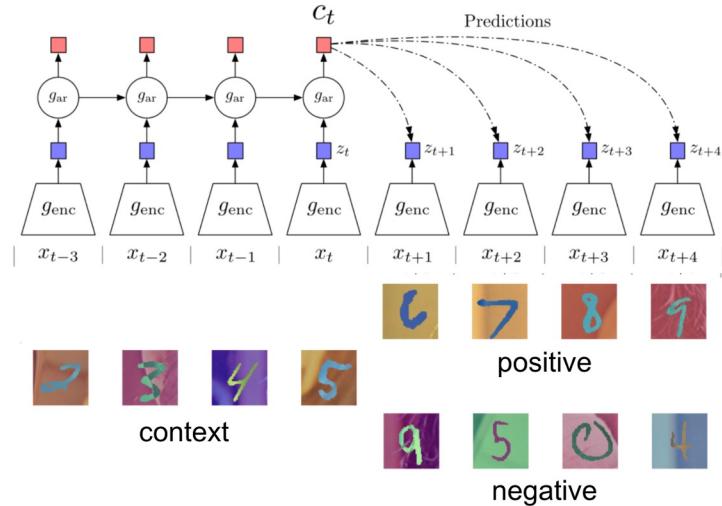
- Decouples negative sample size from minibatch size; allows large batch training without TPU
- MoCo-v2 combines the key ideas from SimCLR, i.e., nonlinear projection head, strong data augmentation, with momentum contrastive learning



# Summary: Contrastive Representation Learning

CPC: sequence-level contrastive learning

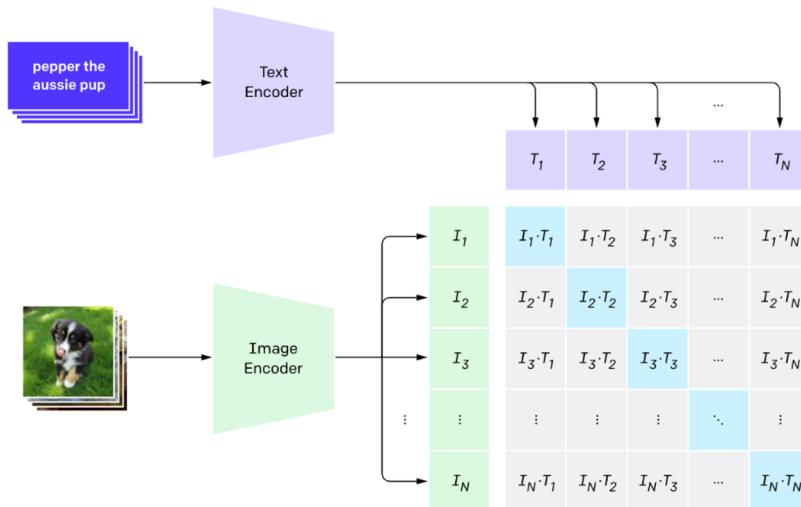
- Contrast “right” sequence with “wrong” sequence.
- InfoNCE loss with a time-dependent score function.
- Can be applied to a variety of learning problems, but not as effective in learning image representations compared to instance-level methods.



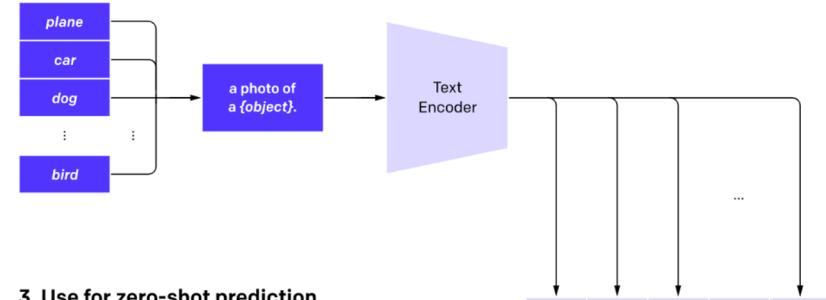
# Other examples

## Contrastive learning between image and natural language sentences

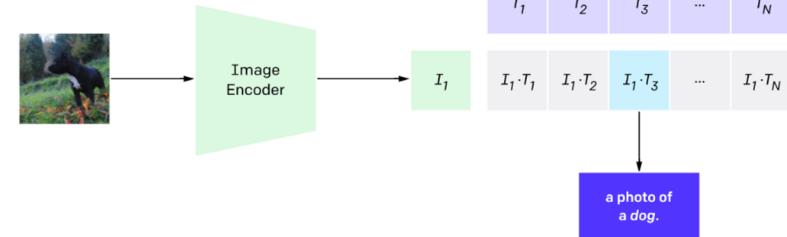
### 1. Contrastive pre-training



### 2. Create dataset classifier from label text



### 3. Use for zero-shot prediction

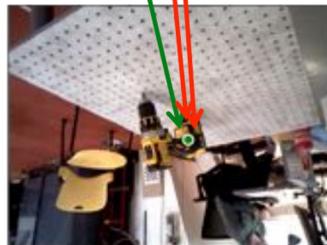
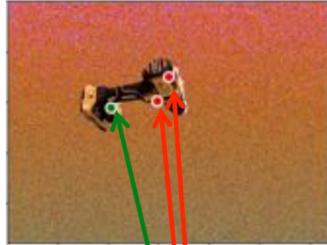


CLIP (*Contrastive Language–Image Pre-training*) Radford *et al.*, 2021

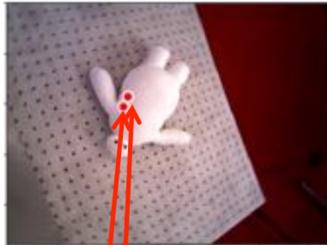
# Other examples

Contrastive learning on pixel-wise feature descriptors

(c) *Background Randomization*



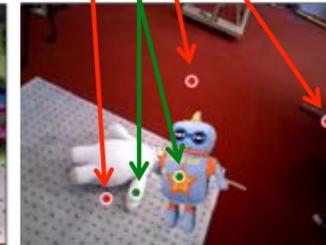
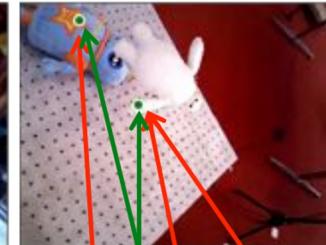
(d) *Cross Object Loss*



(e) *Direct Multi Object*



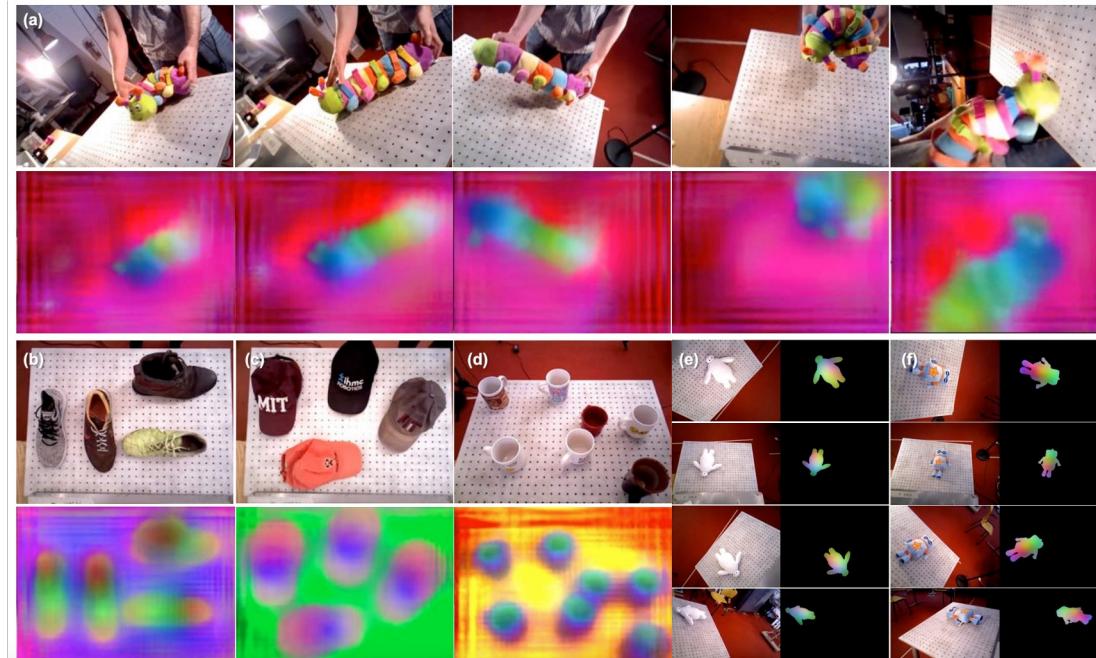
(f) *Synthetic Multi Object*



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Dense Object Net, Florence et al., 2018

# Other examples



Dense Object Net, Florence et al., 2018

**Next Lecture:** Large Vision and Language Model