

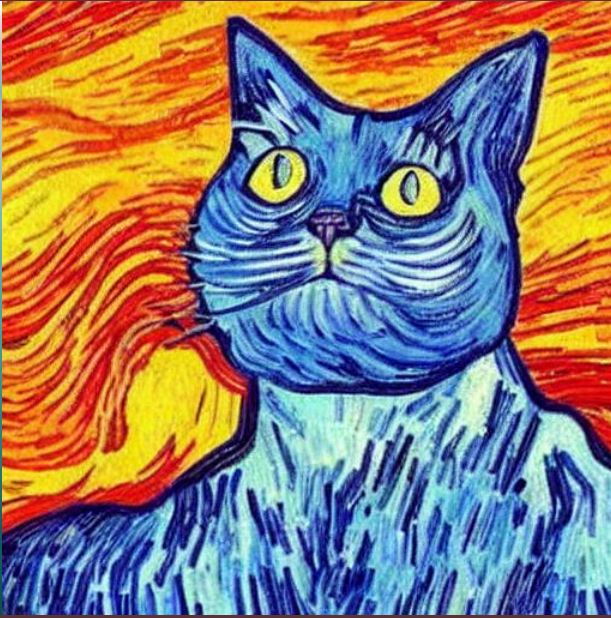
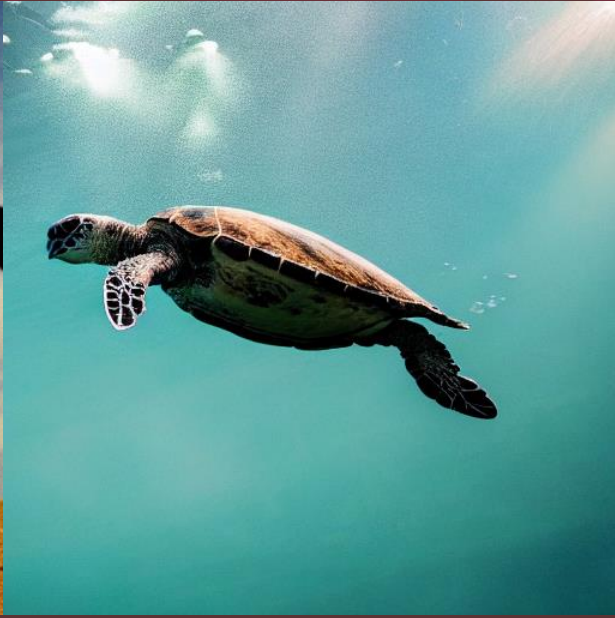
Stable Diffusion

Zhedong Zheng

Credited to Machine Learning from Scratch

Jan 2024





What's the deal with all these pictures?



These pictures were generated by **Stable Diffusion**,
a recent diffusion generative model.

It can turn text prompts (e.g. “an astronaut riding a horse”) into images.
It can also do a variety of other things!

You may have also heard of DALL·E 2, which works in a similar way.



Why should we care?

Could be a model of imagination.

Similar techniques could be used to generate any number of things (e.g. neural data).

"a lovely cat running in the desert in Van Gogh style, trending art."

It's cool!



How does it work?

It's complicated...
but here's the high-level idea.

"Batman eating pizza
in a diner"

What do we need?

"bad stick figure drawing"

Example pictures of people

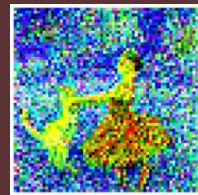


1. Method of learning to generate new stuff given many examples

What do we need?

2. Way to link text and images

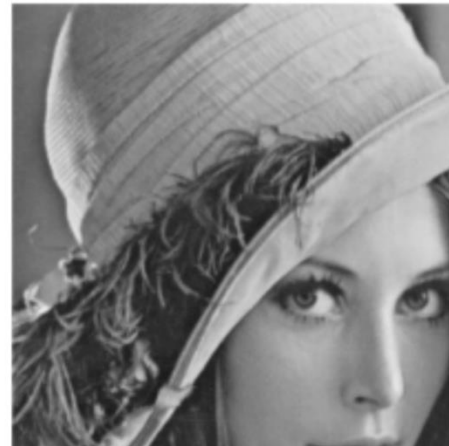
“cool professor person”



$z[0:3, :, :]$

3. Way to compress images (for speed in training and generation)

What do we need?



(a) Original image



(b) Low frequency component

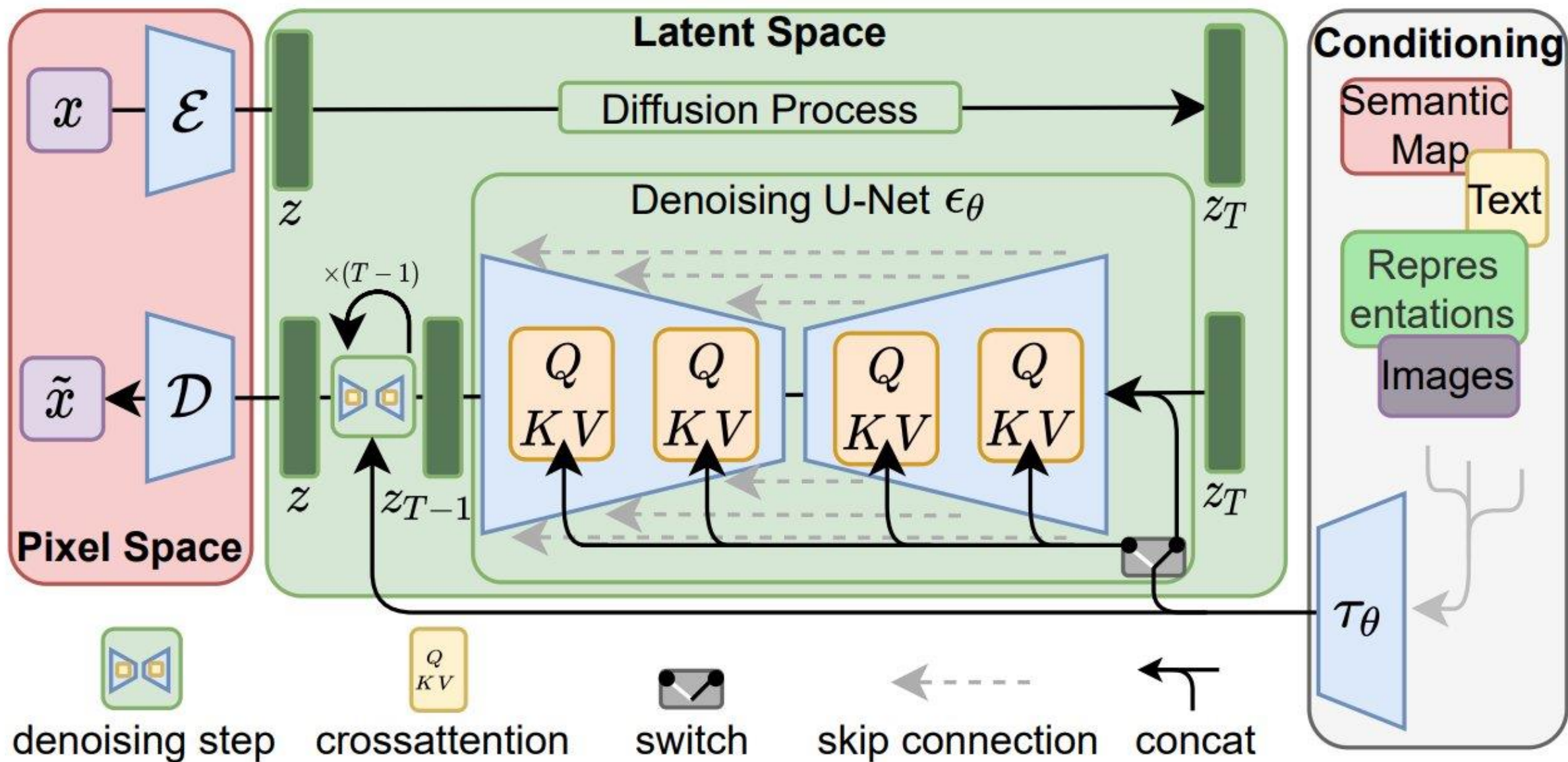
4. Way to add in good image-related inductive biases...

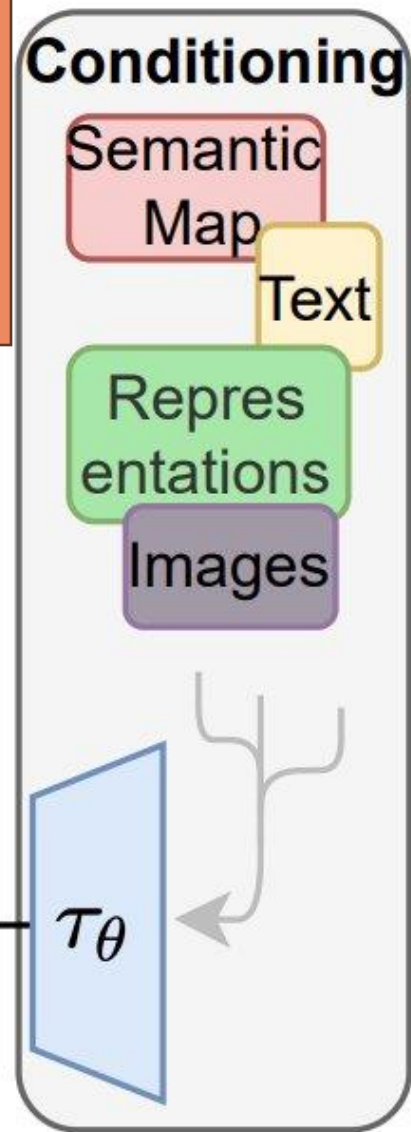
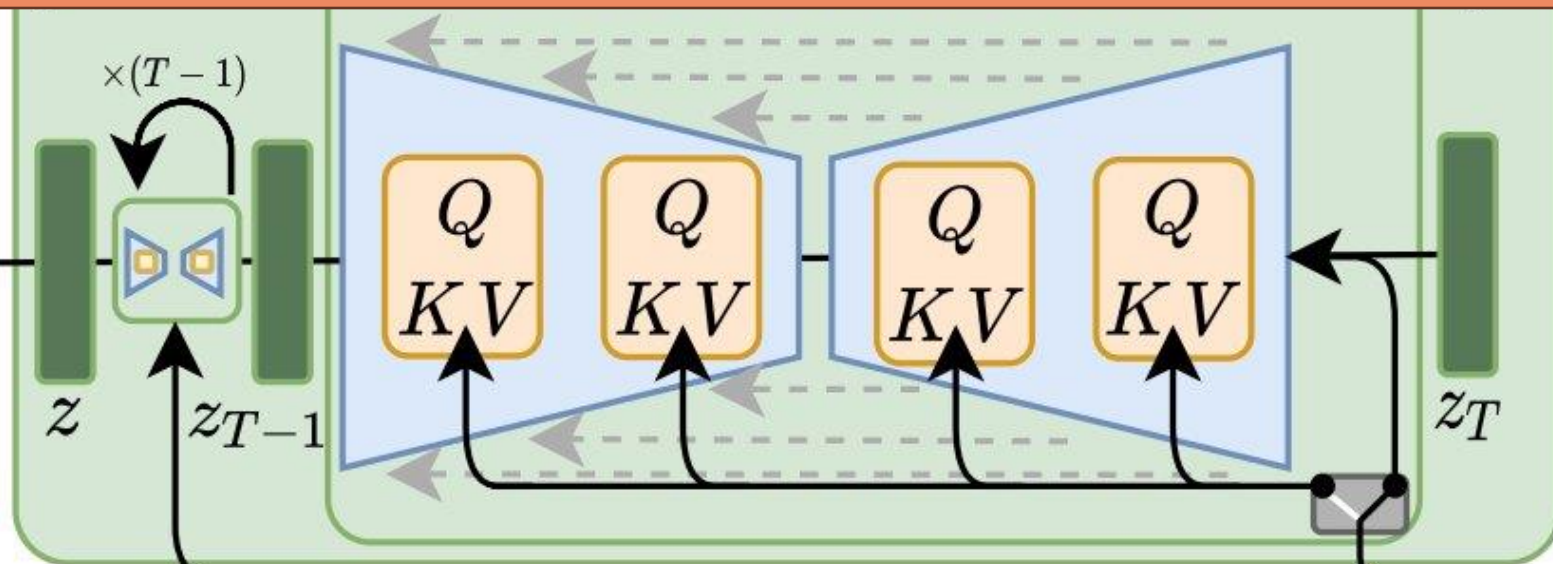
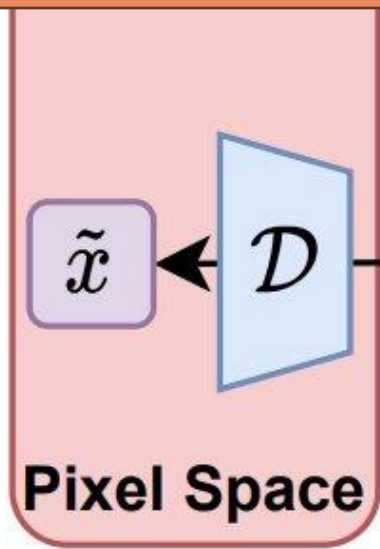
...since when you're generating something new, you need a way to safely go beyond the images you've seen before.

What do we need?

- | | |
|---|---|
| 1. Method of learning to generate new stuff | Forward/reverse diffusion |
| 2. Way to link text and images | Text-image representation model |
| 3. Way to compress images | Autoencoder |
| 4. Way to add in good inductive biases | U-net architecture + 'attention' |

Making a 'good' generative model is about making all these parts work together well!





Stable Diffusion in Action



Cartoon with StableDiffusion + Cartoon



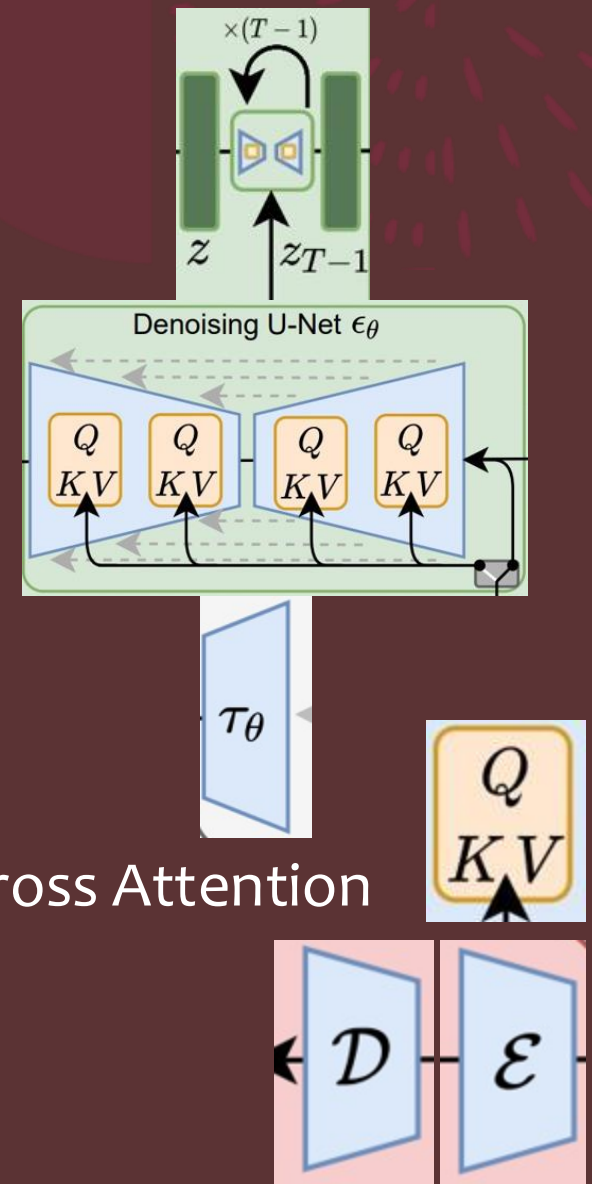
https://www.reddit.com/r/StableDiffusion/comments/xcij7u/sd_img2img_after_effects_i_generated_2_images_and/

Some Resources

- Diffusion model in general
 - [What are Diffusion Models? | Lil'Log](#)
 - [Generative Modeling by Estimating Gradients of the Data Distribution | Yang Song](#)
- Stable diffusion
 - Annotated & simplified code: [U-Net for Stable Diffusion \(labml.ai\)](#)
 - Illustrations: [The Illustrated Stable Diffusion – Jay Alammam](#)
- Attention & Transformers
 - [The Illustrated Transformer](#)

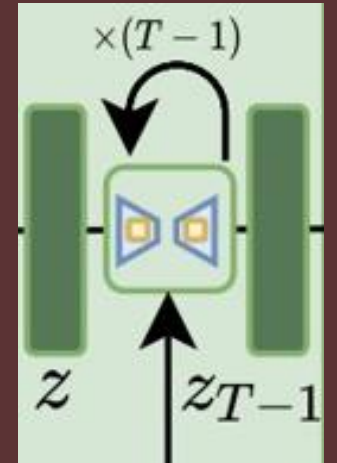
Outline


- Stable Diffusion is cool!
- Build Stable Diffusion “from Scratch”
 - Principle of Diffusion models (sampling, learning)
 - Diffusion for Images – UNet architecture
 - Understanding prompts – Word as vectors, CLIP
 - Let words modulate diffusion – Conditional Diffusion, Cross Attention
 - Diffusion in latent space – AutoEncoderKL
 - Training on Massive Dataset. – LAION 5Billion
- Let’s try ourselves.



Principle of Diffusion Models

Learning to generate by iterative denoising.





*“Creating noise from data is easy;
Creating data from noise is generative modeling.”*

-- Song, Yang

Diffusion models

- Forward diffusion (noising)
 - $x_0 \rightarrow x_1 \rightarrow \cdots x_T$
 - Take a data distribution $x_0 \sim p(x)$, turn it into noise by diffusion $x_T \sim \mathcal{N}(0, \sigma^2 I)$



x_0

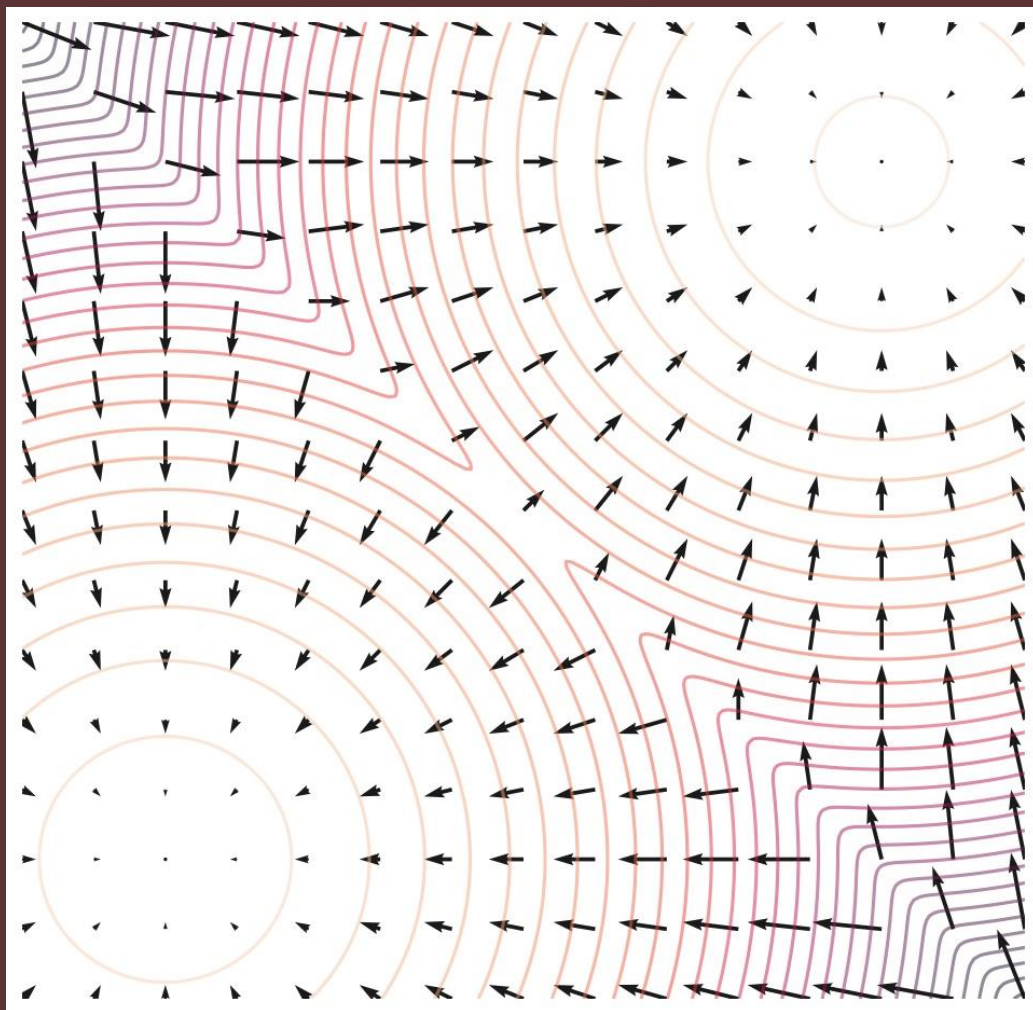
x_1

x_{T-1}

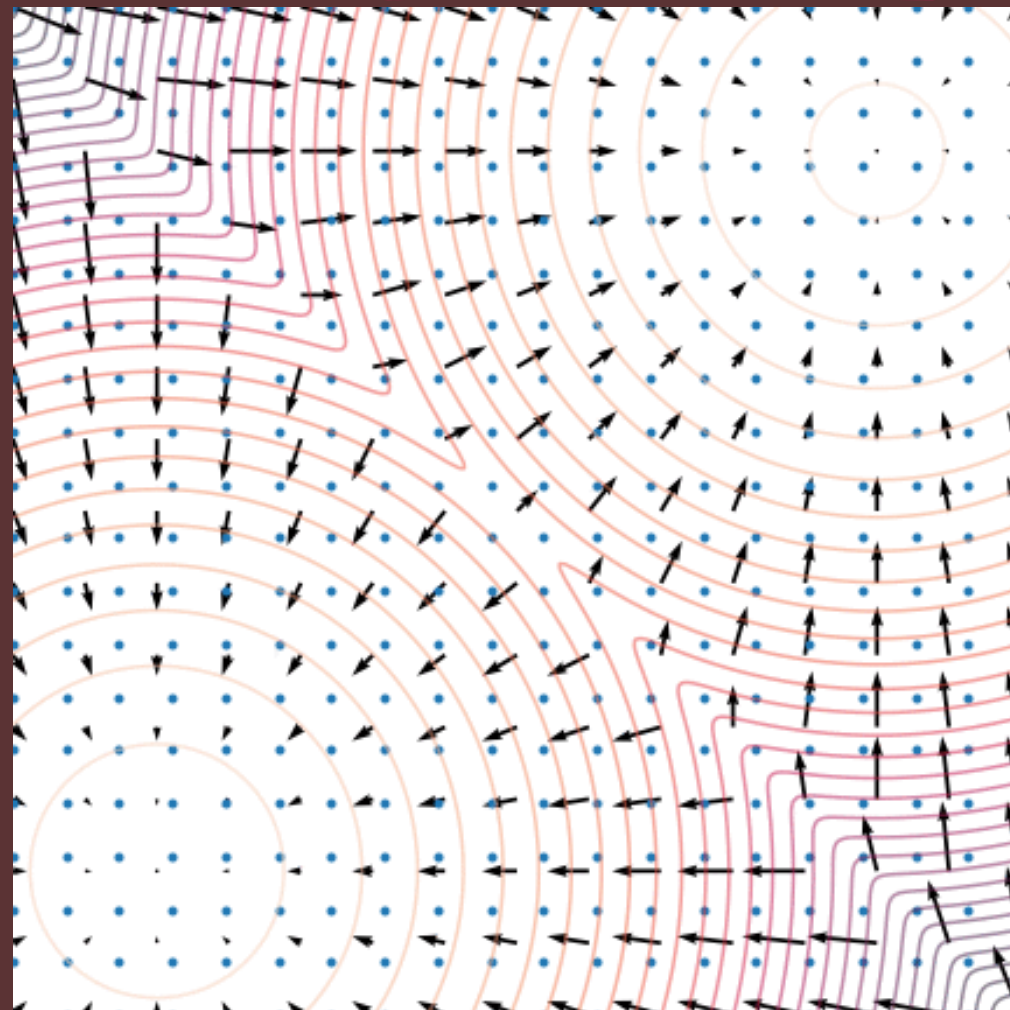
x_T

- Reverse diffusion (denoising)
 - $x_T \rightarrow x_{T-1} \rightarrow \cdots x_0$
 - Sample from the noise distribution $x_T \sim \mathcal{N}(0, \sigma^2 I)$, reverse the diffusion process to generate data $x_0 \sim p(x)$

Animation for the Reverse Diffusion



Score Vector Field



Reverse Diffusion guided by the score vector field

Training diffusion model = Learning to denoise

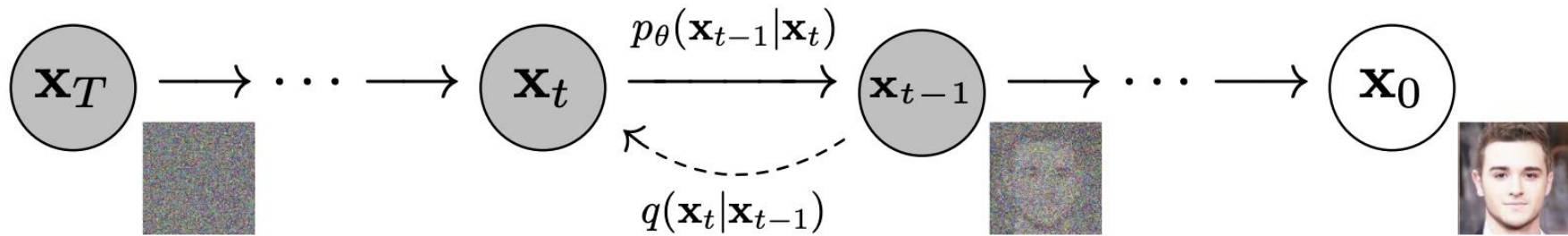


Figure 2: The directed graphical model considered in this work.

$$p(\mathbf{x}_T) = \mathcal{N}(\mathbf{x}_T; \mathbf{0}, \mathbf{I}):$$

$$p_\theta(\mathbf{x}_{t-1}|\mathbf{x}_t) := \mathcal{N}(\mathbf{x}_{t-1}; \boldsymbol{\mu}_\theta(\mathbf{x}_t, t), \boldsymbol{\Sigma}_\theta(\mathbf{x}_t, t))$$

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

Increasing Beta
(0.0001 -> 0.002)

Training diffusion model = Learning to denoise

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
```

Algorithm 2 Sampling

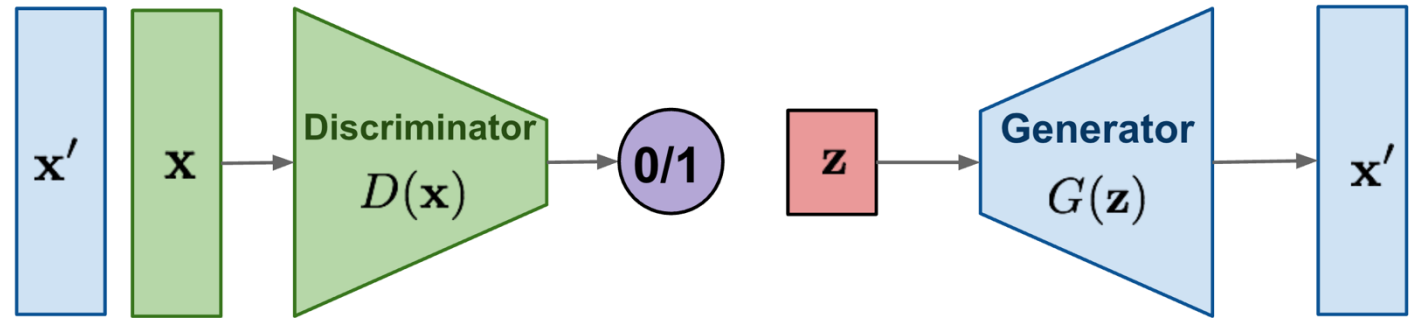
```
1:  $\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$   
2: for  $t = T, \dots, 1$  do  
3:    $\mathbf{z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$  if  $t > 1$ , else  $\mathbf{z} = \mathbf{0}$   
4:    $\mathbf{x}_{t-1} = \frac{1}{\sqrt{\alpha_t}} \left( \mathbf{x}_t - \frac{1 - \alpha_t}{\sqrt{1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

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

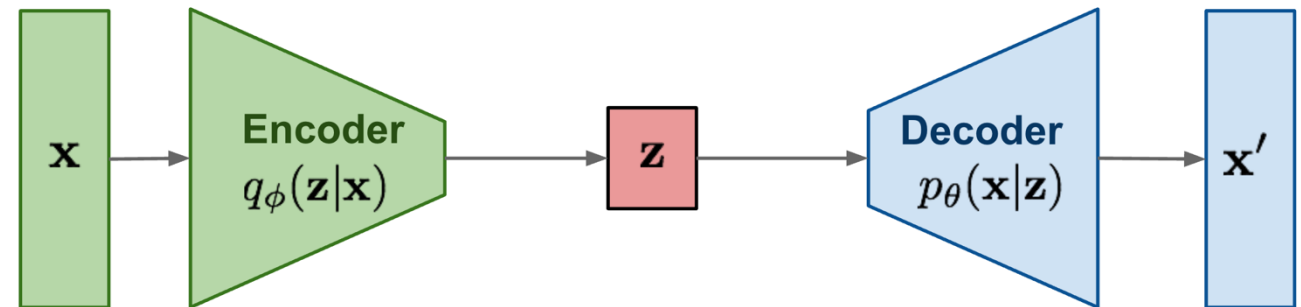
$$\alpha_t := 1 - \beta_t \text{ and } \bar{\alpha}_t := \prod_{s=1}^t \alpha_s,$$

Comparing Generative Models

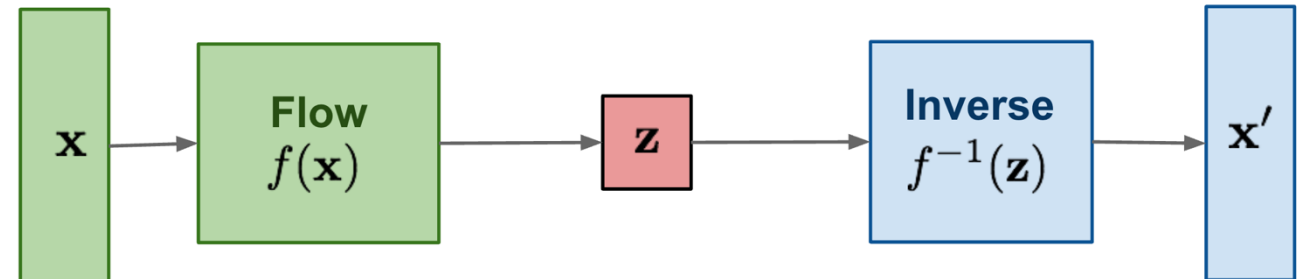
GAN: Adversarial training



VAE: maximize variational lower bound



Flow-based models:
Invertible transform of distributions



Diffusion models:
Gradually add Gaussian noise and then reverse



Diffusion vs GAN / VAE

GAN

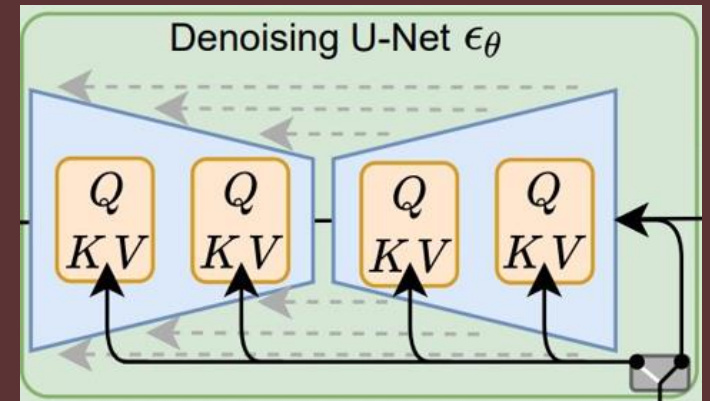
- One shot generation. Fast.
- Harder to control in one pass.
- Adversarial min-max objective. Can collapse.

Diffusion

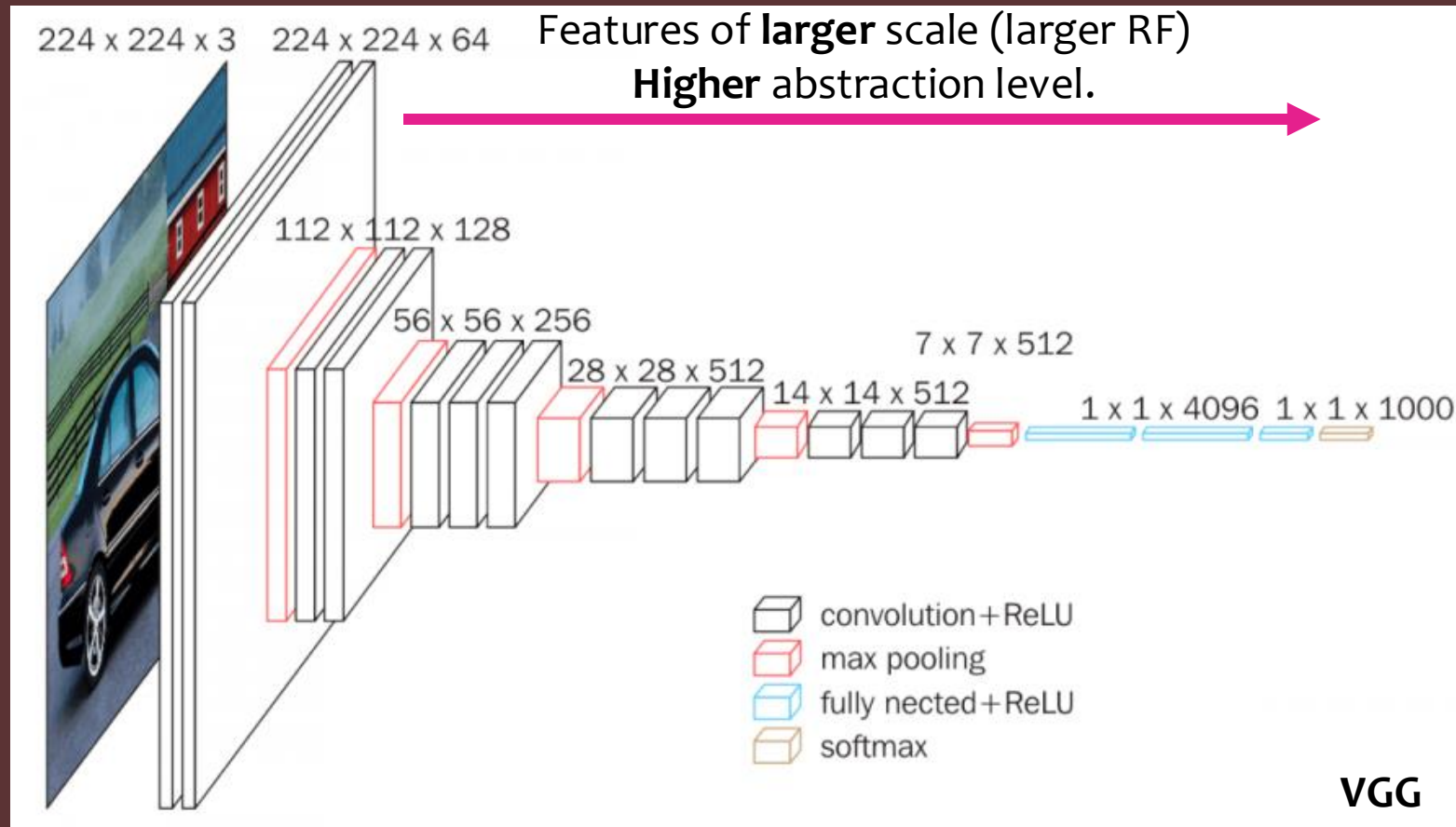
- Multi-iteration generation. Slow.
- Easier to control during generation.
- Simple objective, no adversary in training.

Modelling Score function over Image Domain

Introducing UNet



Convolutional Neural Network



- CNN parametrizes function over images
- Motivation
 - Features are translational invariant
 - Extract feature at different scale / abstraction level
- Key modules
 - Convolution
 - Downsampling (Max-pool)

CNN + inverted CNN \Rightarrow UNet

CNN + inverted CNN \Rightarrow UNet

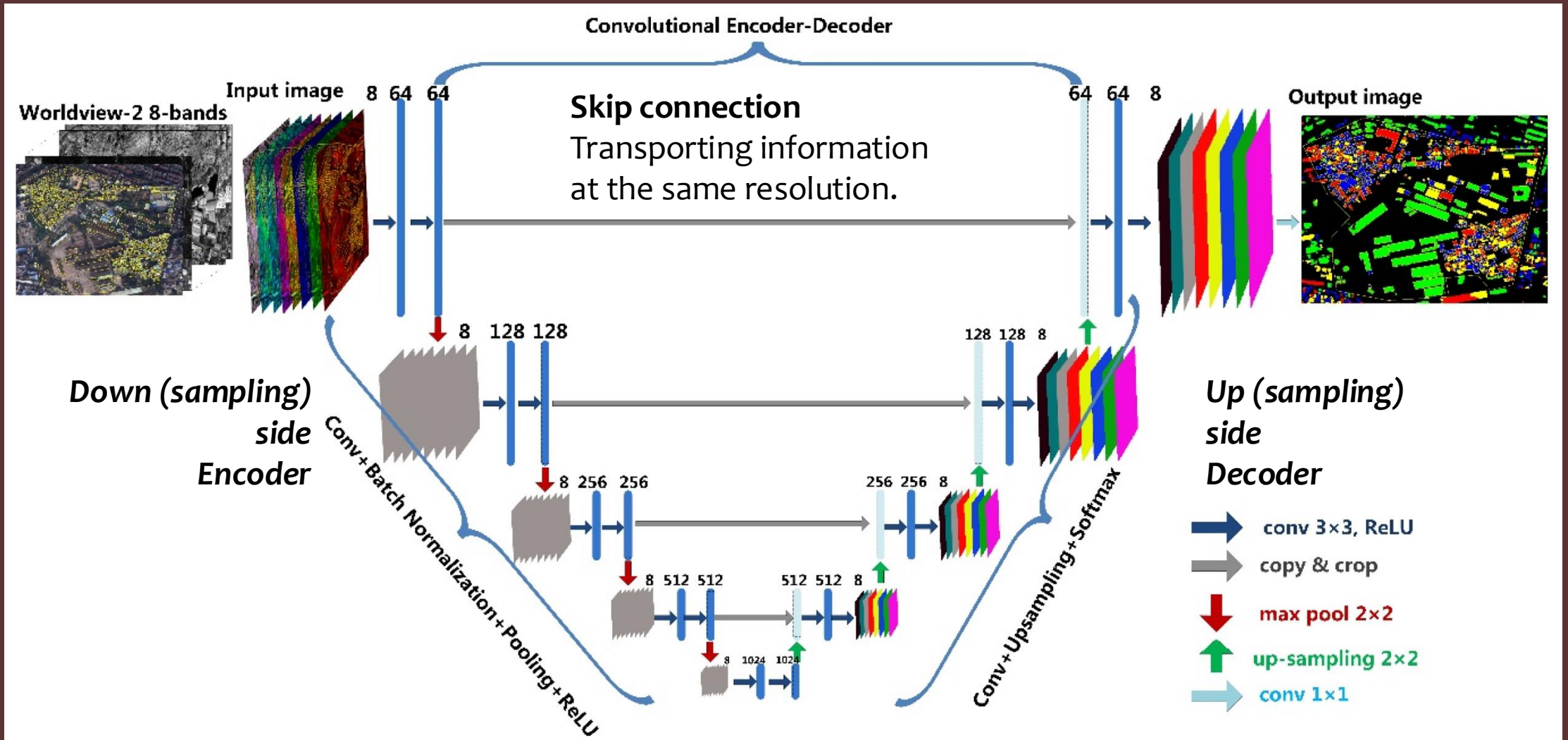
Down Sampling
Convolution

CNN + inverted CNN \Rightarrow UNet

Up Sampling
Transposed Convolution

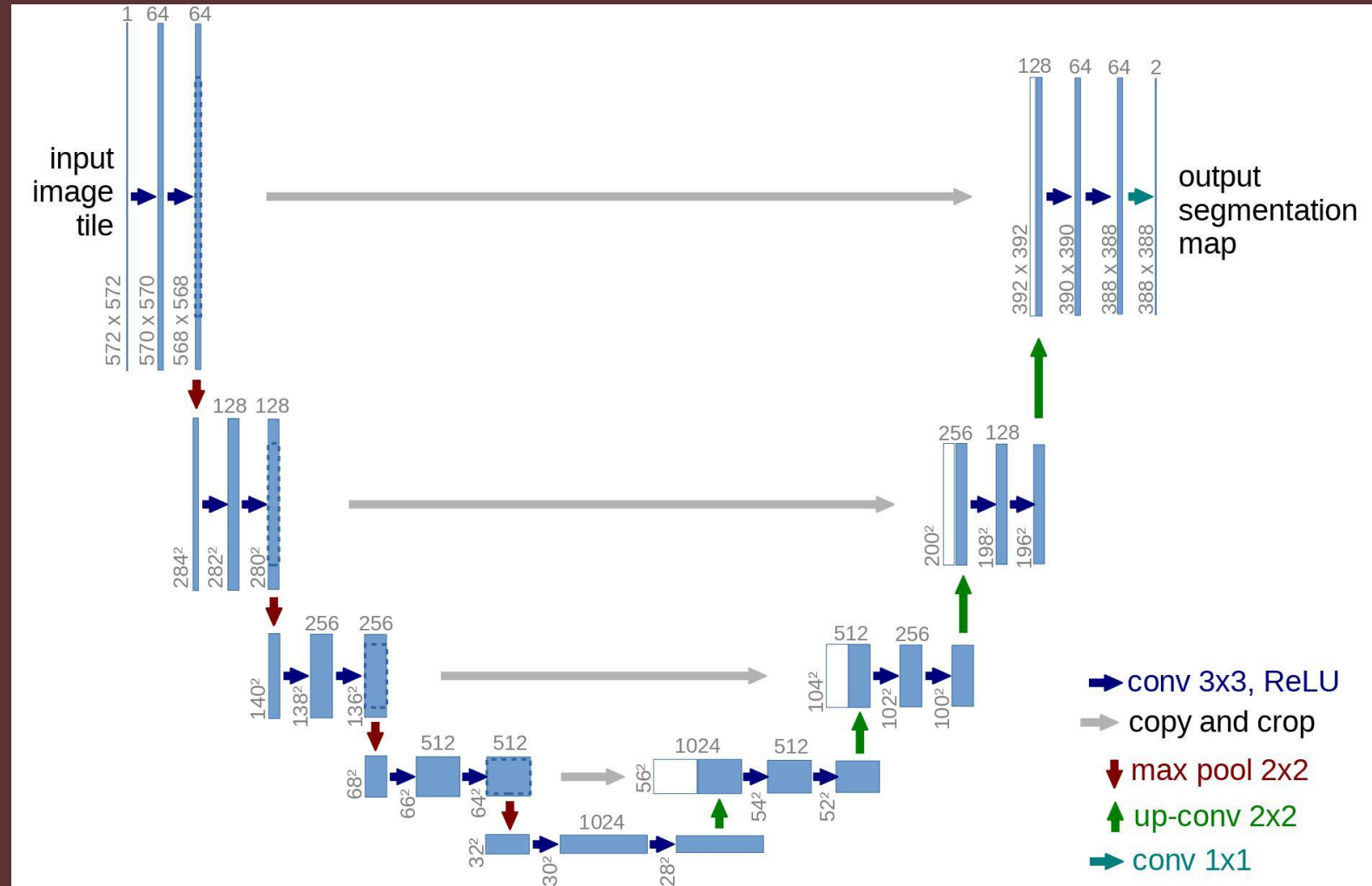
- Inverted CNN (generator) can generate images.
- CNN + inverted CNN could model Image \rightarrow Image function.

UNet: a natural architecture for image-to-image function



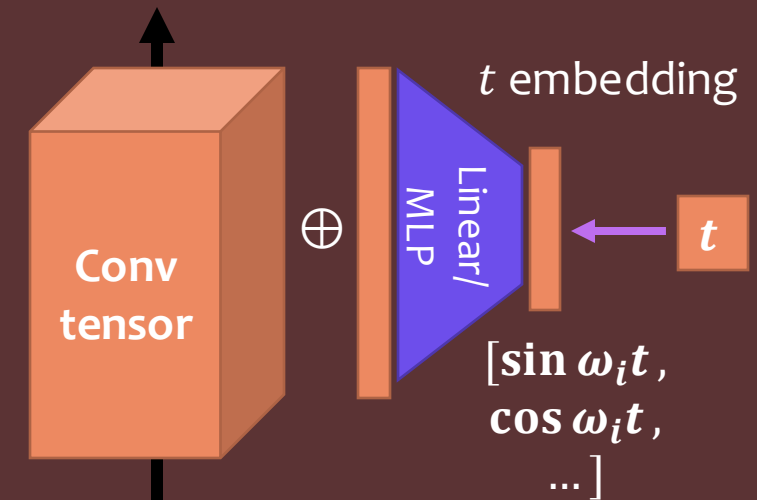
Key Ingredients of UNet

- Convolution operation
 - Save parameter, spatial invariant
- Down/Up sampling
 - Multiscale / Hierarchy
 - Learn modulation at multi scale and multi-abstraction levels.
- Skip connection
 - **No bottleneck**
 - Route feature of the same scaledirectly.
 - Cf. AutoEncoder has bottleneck

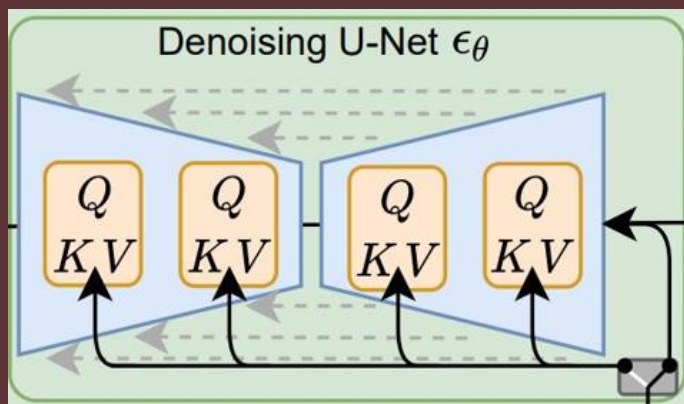


Note: Add Time Dependency

- The score function is *time-dependent*.
 - Target: $s(x, t) = \nabla_x \log p(x, t)$
- Add time dependency
 - Assume time dependency is spatially homogeneous.
 - Add one scalar value per channel $f(t)$
 - Parametrize $f(t)$ by MLP / linear of Fourier basis.



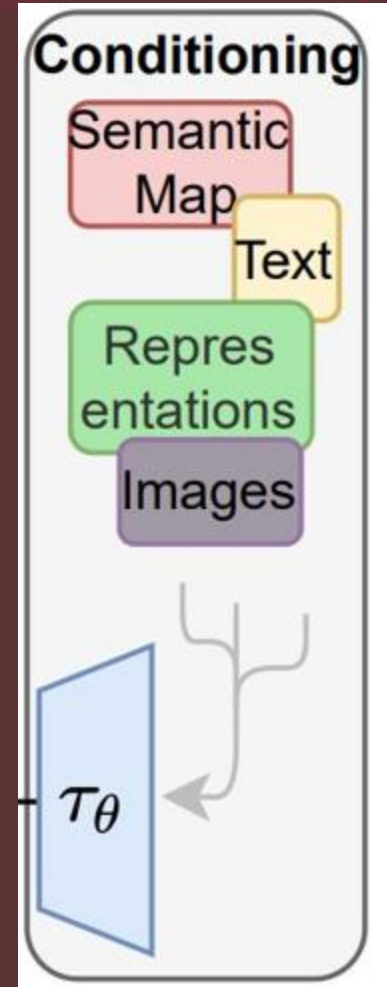
Unet in Stable Diffusion



```
(conv_in): Conv2d(4, 320, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
(time_proj): Timesteps()
(time_embedding): TimestepEmbedding
  (linear_1): Linear(in_features=320, out_features=1280, bias=True)
  (act): SiLU()
  (linear_2): Linear(in_features=1280, out_features=1280, bias=True)
(down_blocks):
  (0): CrossAttnDownBlock2D
  (1): CrossAttnDownBlock2D
  (2): CrossAttnDownBlock2D
  (3): DownBlock2D
(up_blocks):
  (0): UpBlock2D
  (1): CrossAttnUpBlock2D
  (2): CrossAttnUpBlock2D
  (3): CrossAttnUpBlock2D
(mid_block): UNetMidBlock2DCrossAttn
  (attentions):
  (resnets):
(conv_norm_out): GroupNorm(32, 320, eps=1e-05, affine=True)
(conv_act): SiLU()
(conv_out): Conv2d(320, 4, kernel_size=(3, 3), stride=(1, 1), padding=(1, 1))
```

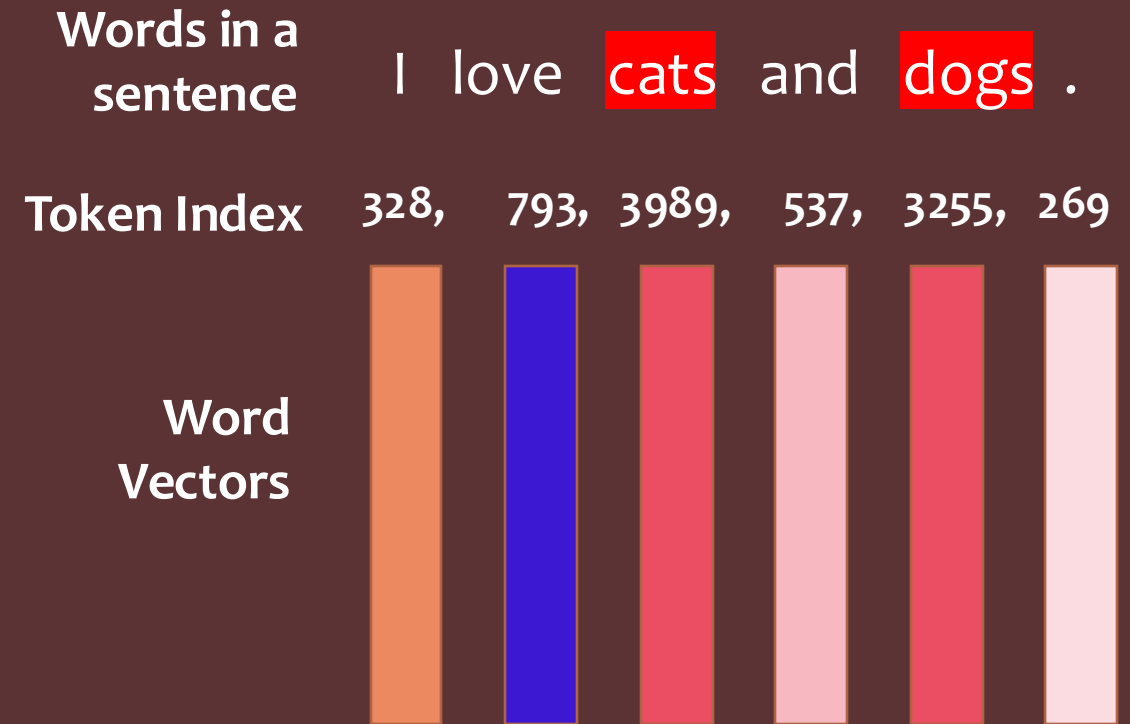
How to understand prompts?

Language / Multimodal Transformer, CLIP!



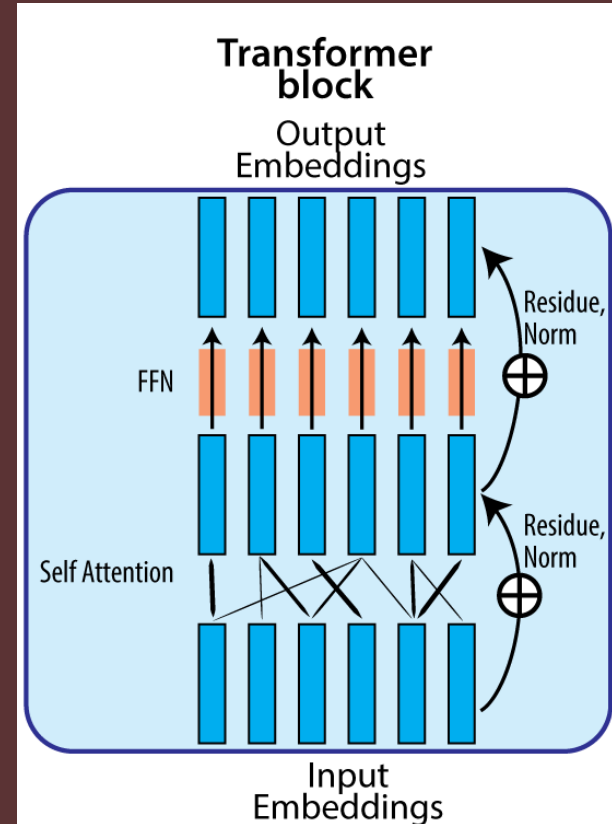
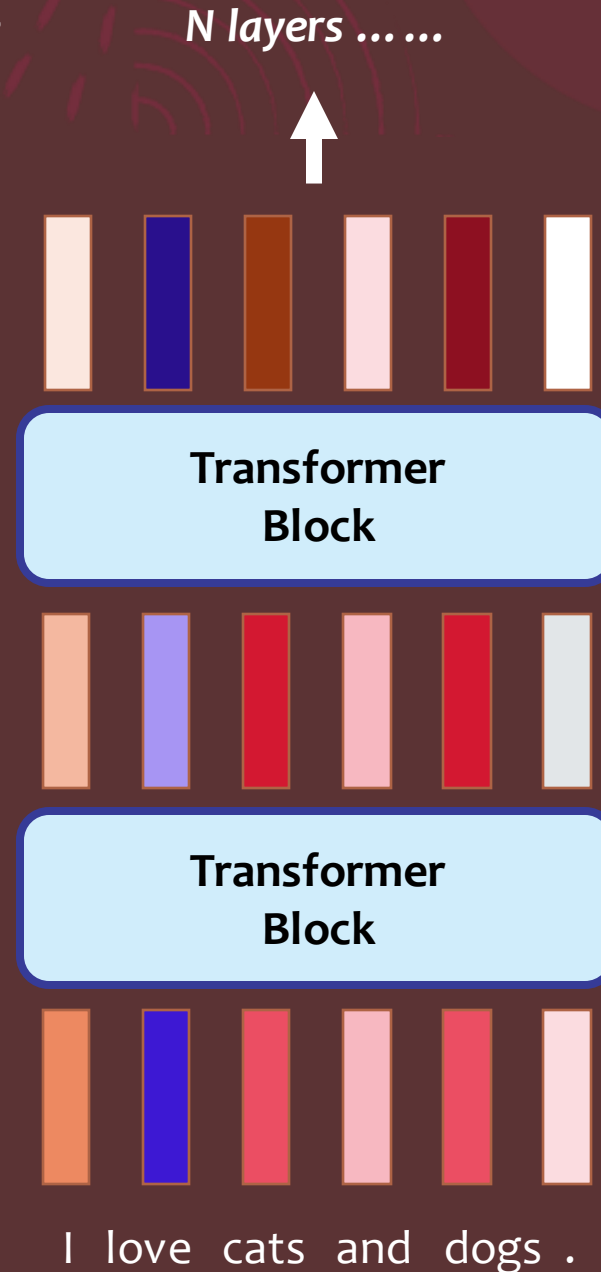
Word as Vectors: Language Model 101

- Unlike pixel, meaning of word are not explicitly in the characters.
- Word can be represented as index in dictionary
 - But index is also meaning less.
- Represent words in a vector space
 - Vector geometry => semantic relation.



Word Vector in Context: RNN / Transformers

- Meaning of word depends on context, not always the same.
 - “I **book** a ticket to buy that **book**.”
- Word vectors should depend on context.
- Transformers let each word “absorb” influence from other words to be “contextualized”

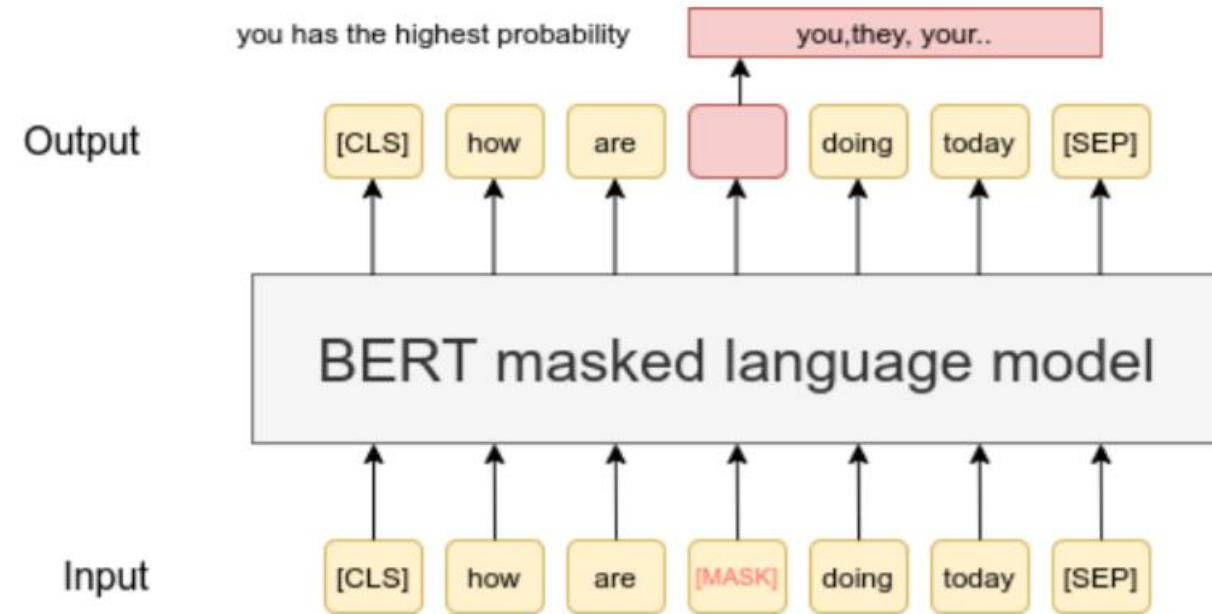


More on attention later...

Learning Word Vectors: GPT & BERT & CLIP

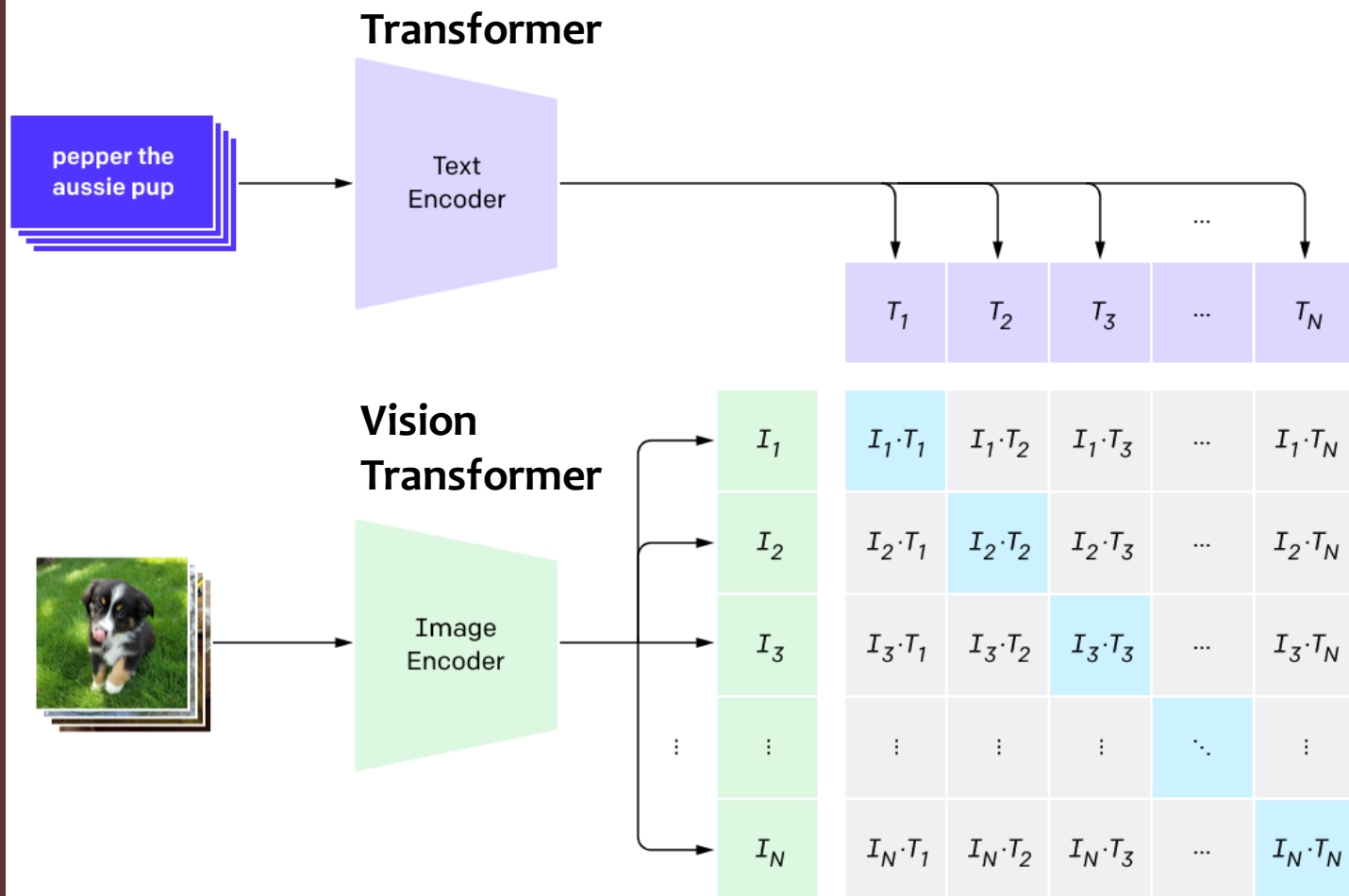
- Self-supervised learning of word representation
 - Predicting missing / next words in a sentence. (BERT, GPT)
 - Contrastive Learning, matching image and text. (CLIP)

Downstream Classifier can decode:
Part of speech, Sentiment, ...



Joint Representation for Vision and Language : CLIP

1. Contrastive pre-training

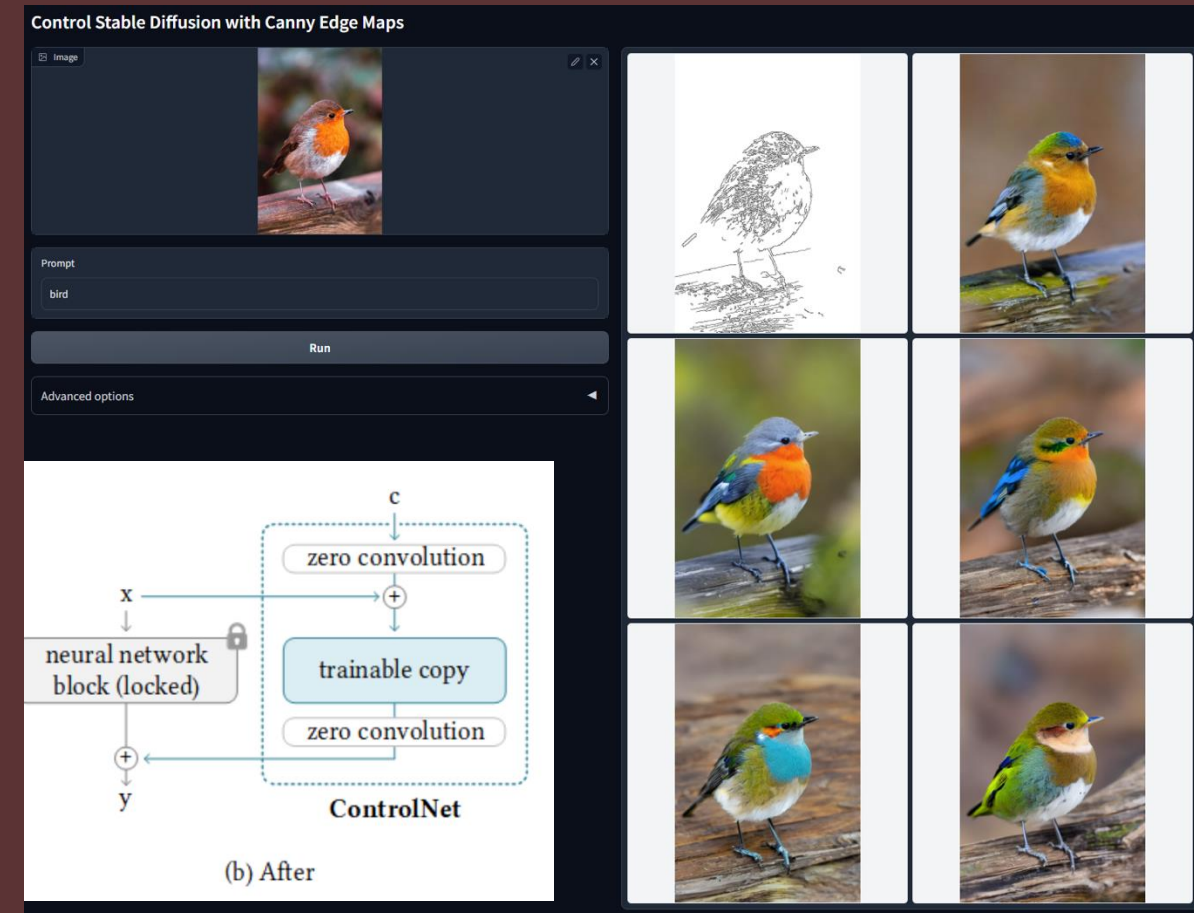
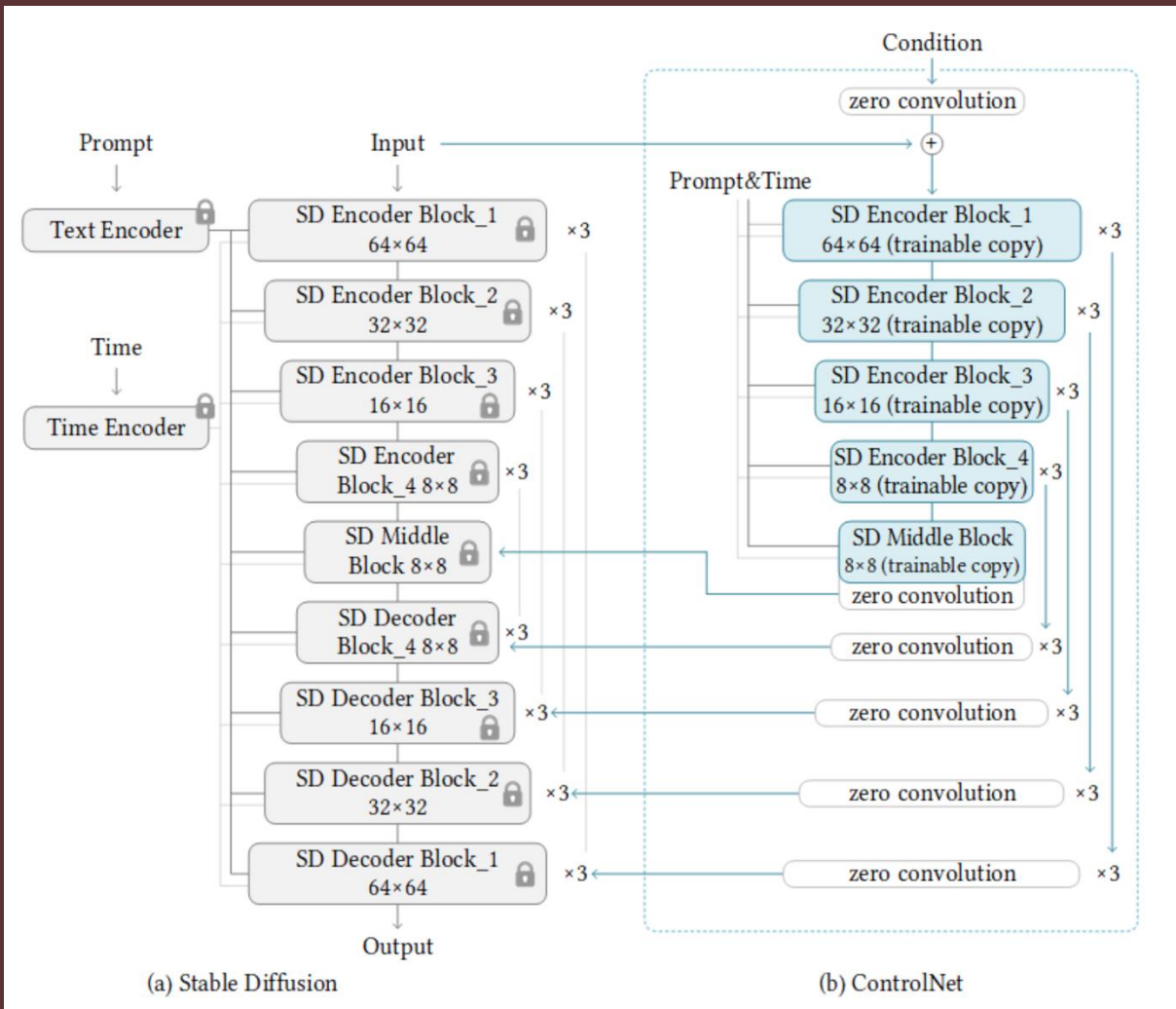


- Learn a joint encoding space for text caption and image
- Maximize representation similarity between an image and its caption.
- Minimize other pairs

Choice of text encoding

- Encoder in Stable Diffusion: pre-trained CLIP ViT-L/14 text encoder
- Word vector can be randomly initialized and learned online.
- Representing other conditional signals
 - Object categories (e.g. Shark, Trout, etc.):
 - 1 vector per class
 - Face attributes (e.g. {female, blonde hair, with glasses, ... }, {male, short hair, dark skin}):
 - set of vectors, 1 vector per attributes
- Time to be creative!!

Conditional Diffusion Models



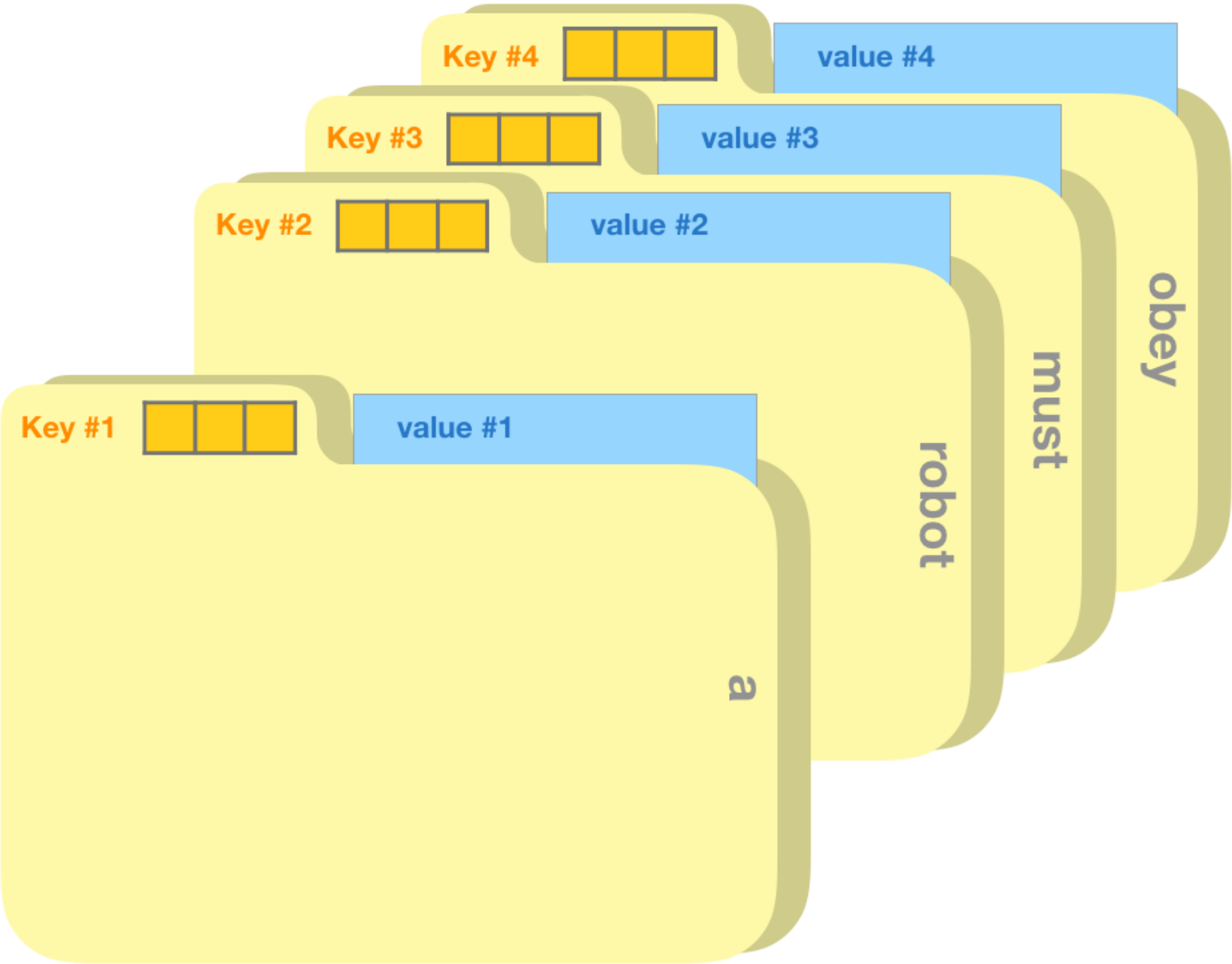
<https://github.com/lllyasviel/ControlNet>

How does text affect diffusion?

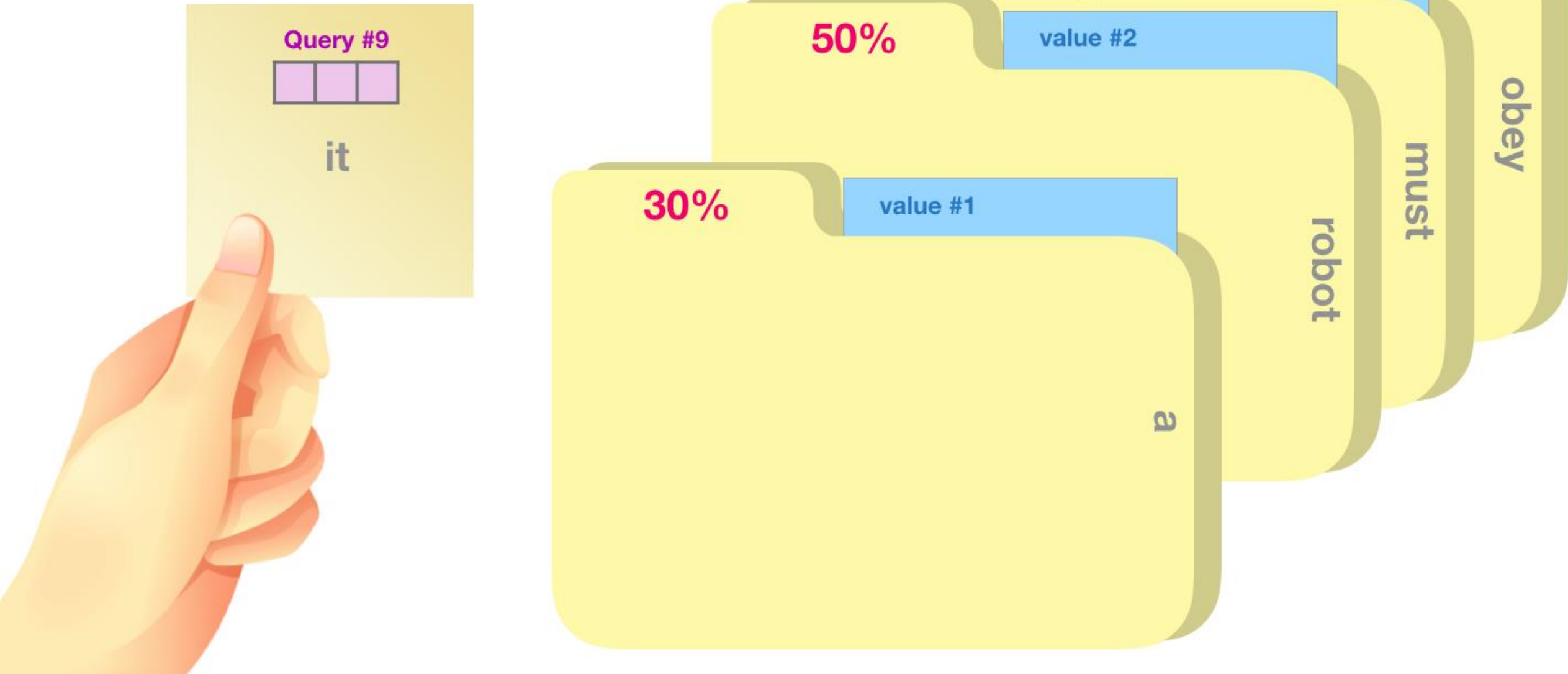
Incoming Cross Attention



“A robot must obey the order given *it*.”



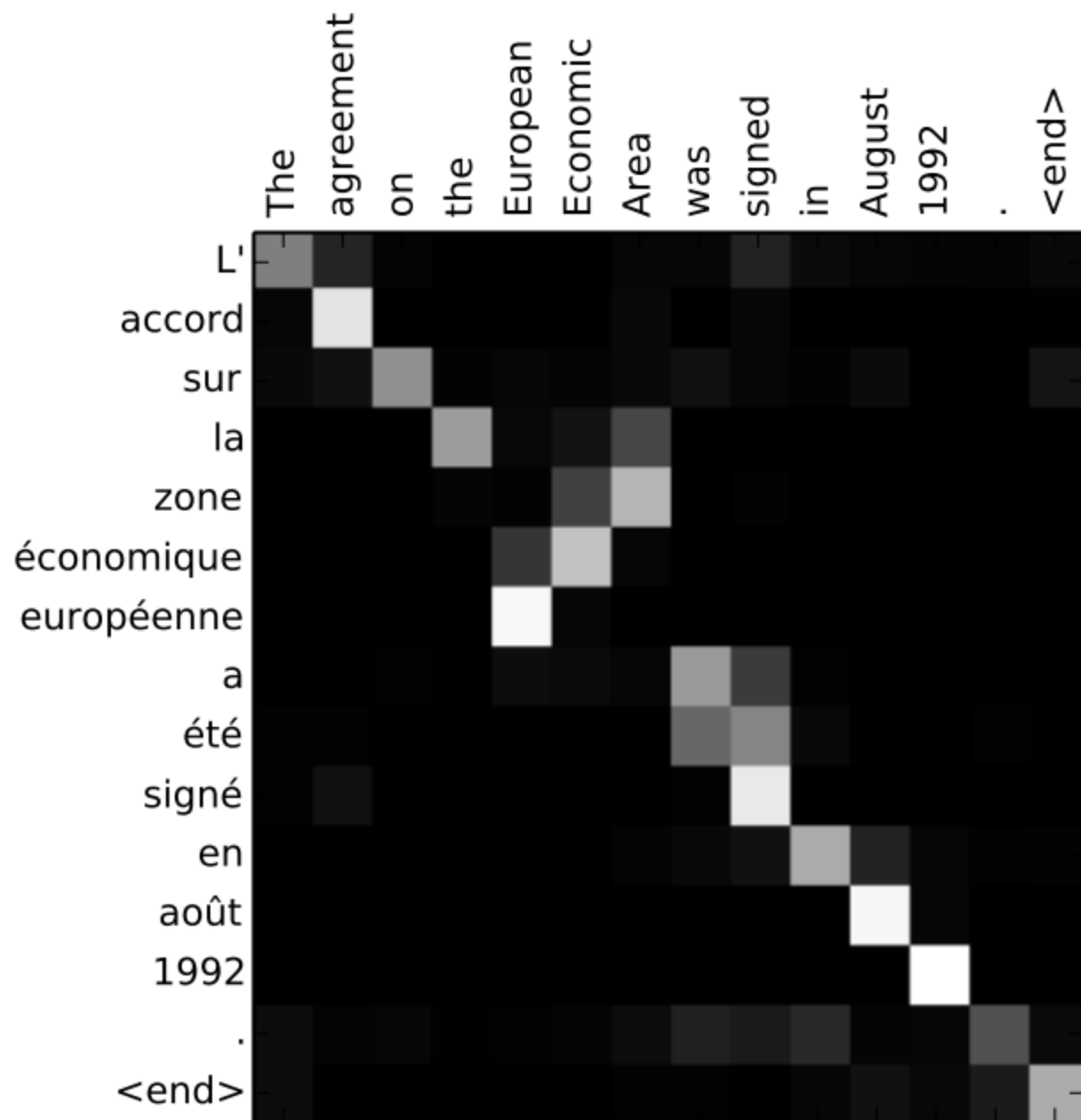
“A robot must obey the order given **it**.”



Visualizing Attention matrix a_{ij}

- French 2 English
- “Learnt to pay Attention”
 - “la zone economique europeenne” -> “the European Economic Area”
 - “a ete signe” -> “was signed”

Attention + RNN



Text2Image as translation

Source language: Words

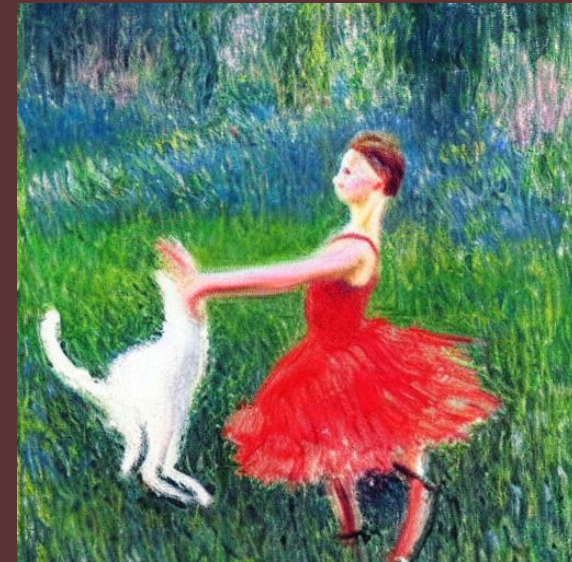
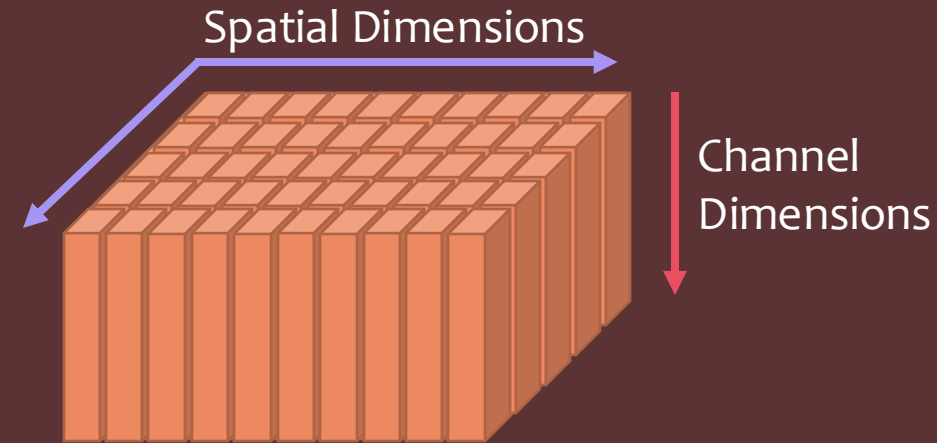
Target language: Images



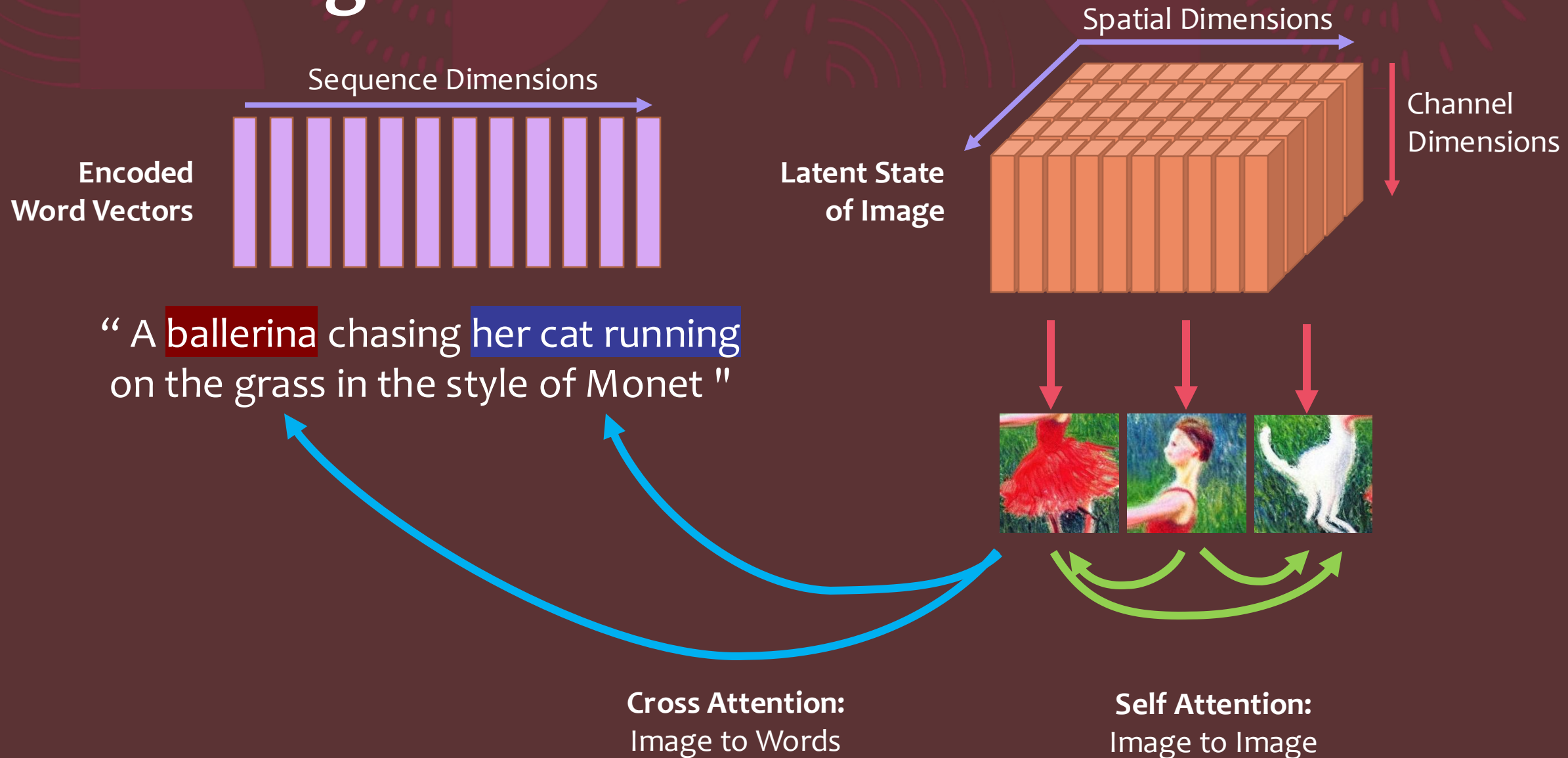
“ A ballerina chasing her cat running on the grass in the style of Monet ”

Latent State of Image

Patch Vectors!

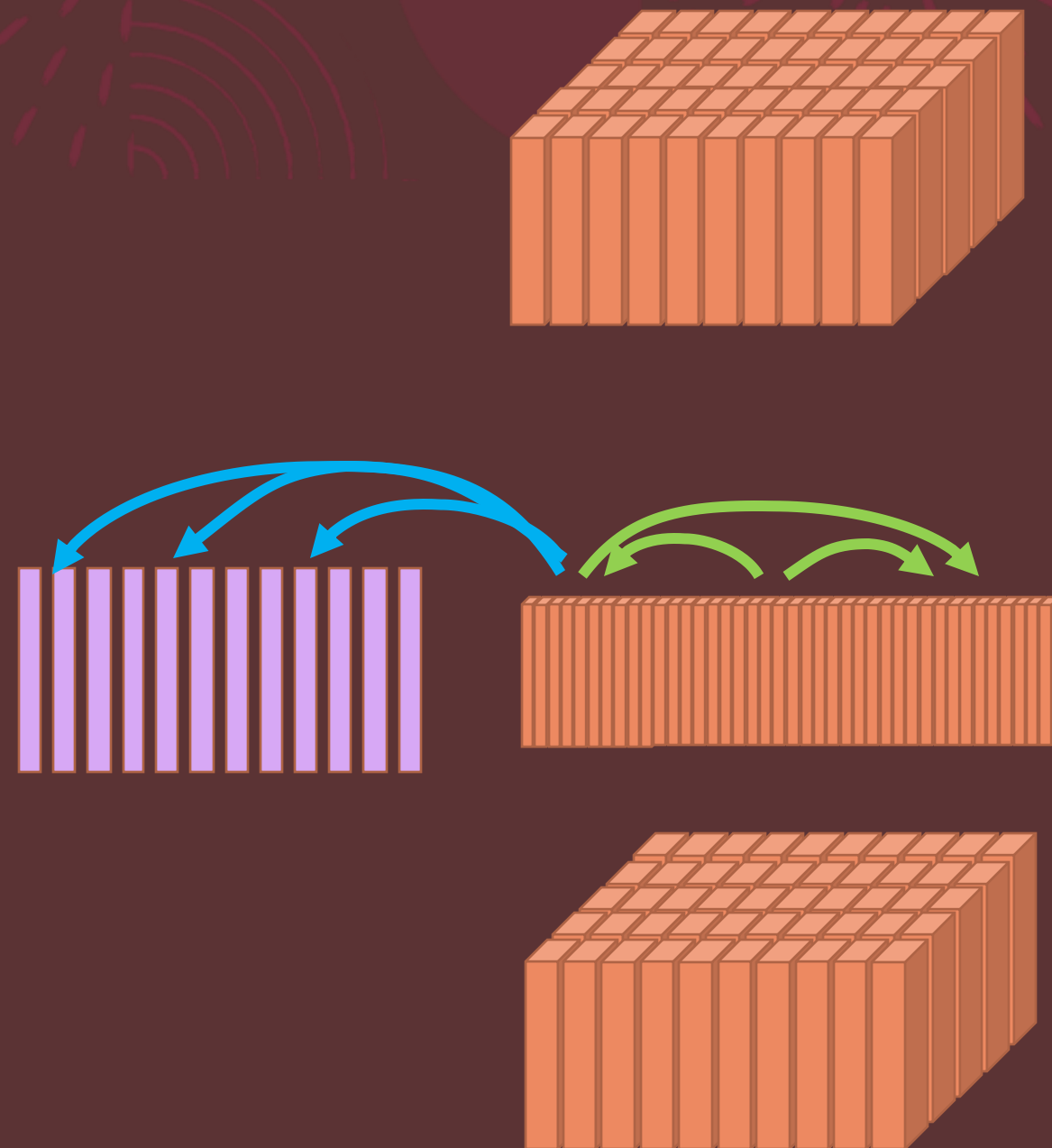


Text2Image as translation



Spatial Transformer

- Rearrange spatial tensor to sequence.
- Cross Attention
- Self Attention
- FFN
- Rearrange back to spatial tensor (same shape)



UNet = Giant Sandwich of Spatial transformer + ResBlock (Conv layer)

- Down blocks

ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

DownSample

ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

DownSample

ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

DownSample

ResBlock

ResBlock

- Up blocks

ResBlock

ResBlock

ResBlock

UpSample

ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

UpSample

ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

UpSample

ResBlock

SpatialTransformer

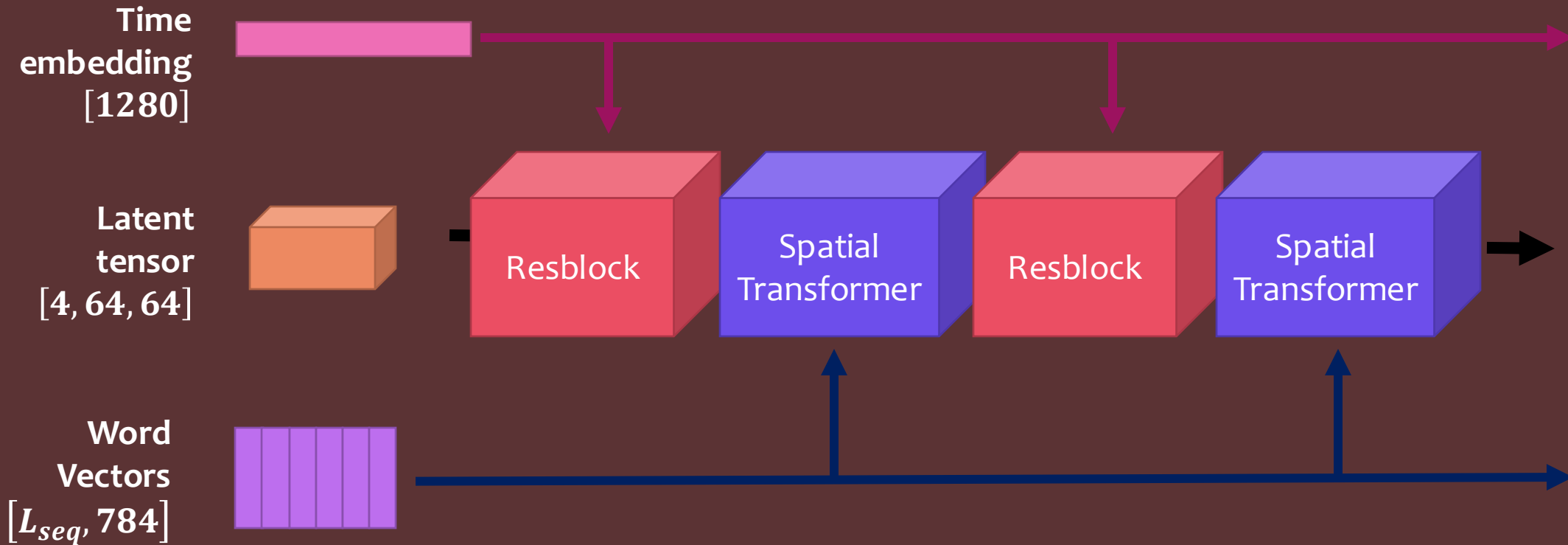
ResBlock

SpatialTransformer

ResBlock

SpatialTransformer

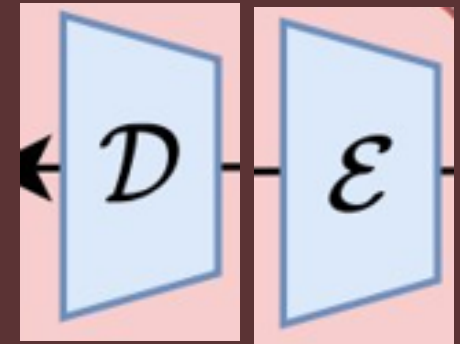
Spatial transformer + ResBlock (Conv layer)



- Alternating Time and Word Modulation
- Alternating Local and Nonlocal operation

Diffusion in Latent Space

Adding in AutoEncoder



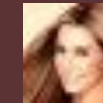
Rombach, R., Blattmann, A., Lorenz, D., Esser, P., & Ommer, B. (2022). High-resolution image synthesis with latent diffusion models, *CVPR*

Diffusion in latent space

- Motivation:
 - Natural images are high dimensional
 - but have many redundant details that could be compressed / statistically filled out
- Division of labor
 - Diffusion model -> Generate low resolution sketch
 - AutoEncoder -> Fill out high resolution details
- Train a VAE model to compress images into latent space.
 - $x \rightarrow z \rightarrow \hat{x}$
- Train diffusion models in latent space of z .

DownSampling

32 pix



$d = 2352$

180 pix

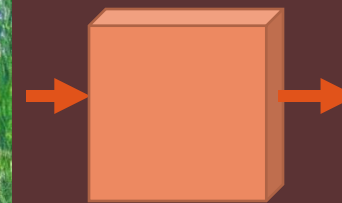


$d = 97200$



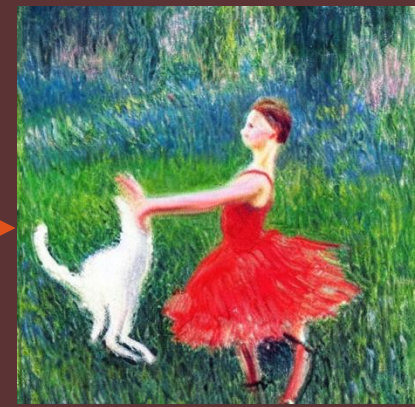
x

$[3,512,512]$



z

$[4,512/f, 512/f]$

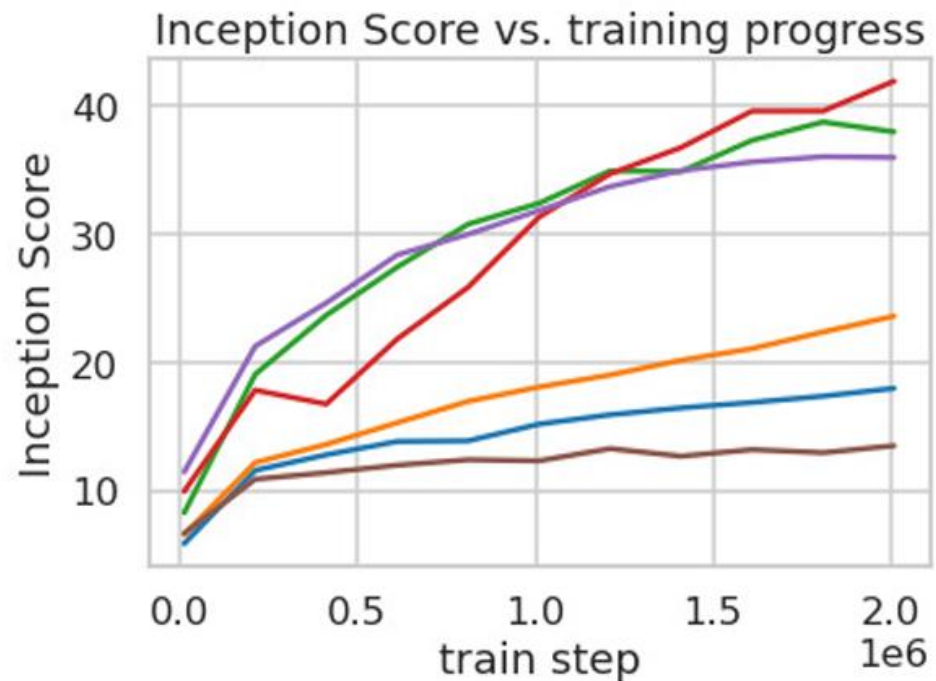
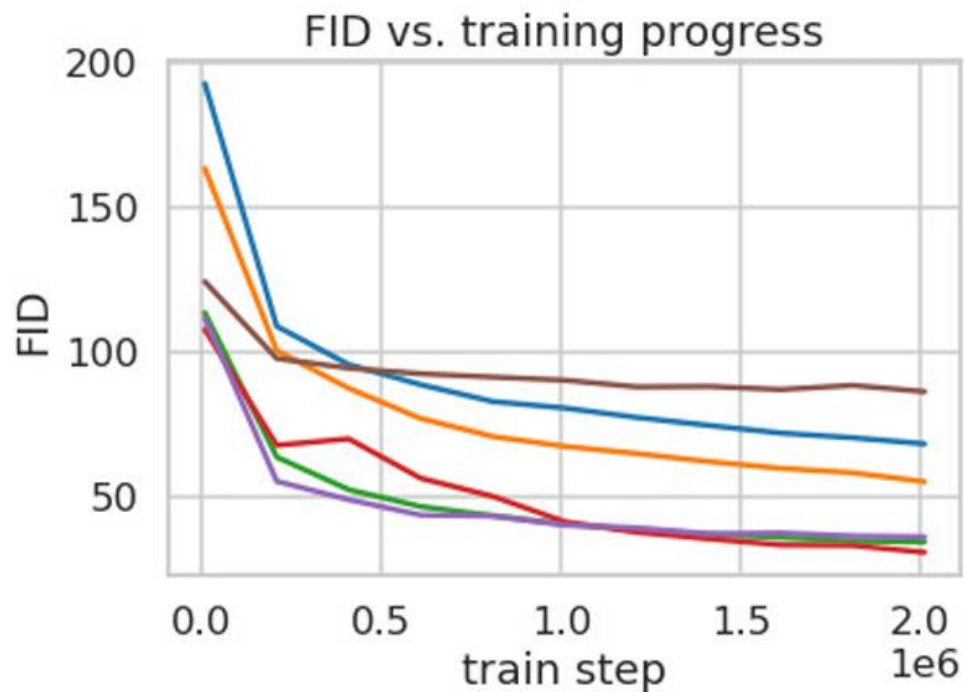


\hat{x}

$[3,512,512]$

Spatial Compression Tradeoff

- LDM- $\{f\}$. f = Spatial downsampling factor
 - Too little compression $f = 1, 2$ or too much compression $f = 32$, makes diffusion hard to train.



Details in Stable Diffusion

- In stable diffusion, spatial downsampling $f = 8$
 - x is (3, 512, 512) image tensor
 - z is (4, 64, 64) latent tensor

Regularizing the Latent Space

- KL regularizer
 - Similar to VAE, make latent distribution like Gaussian distribution.
- VQ regularizer
 - Make the latent representation quantized to be a set of discrete tokens.



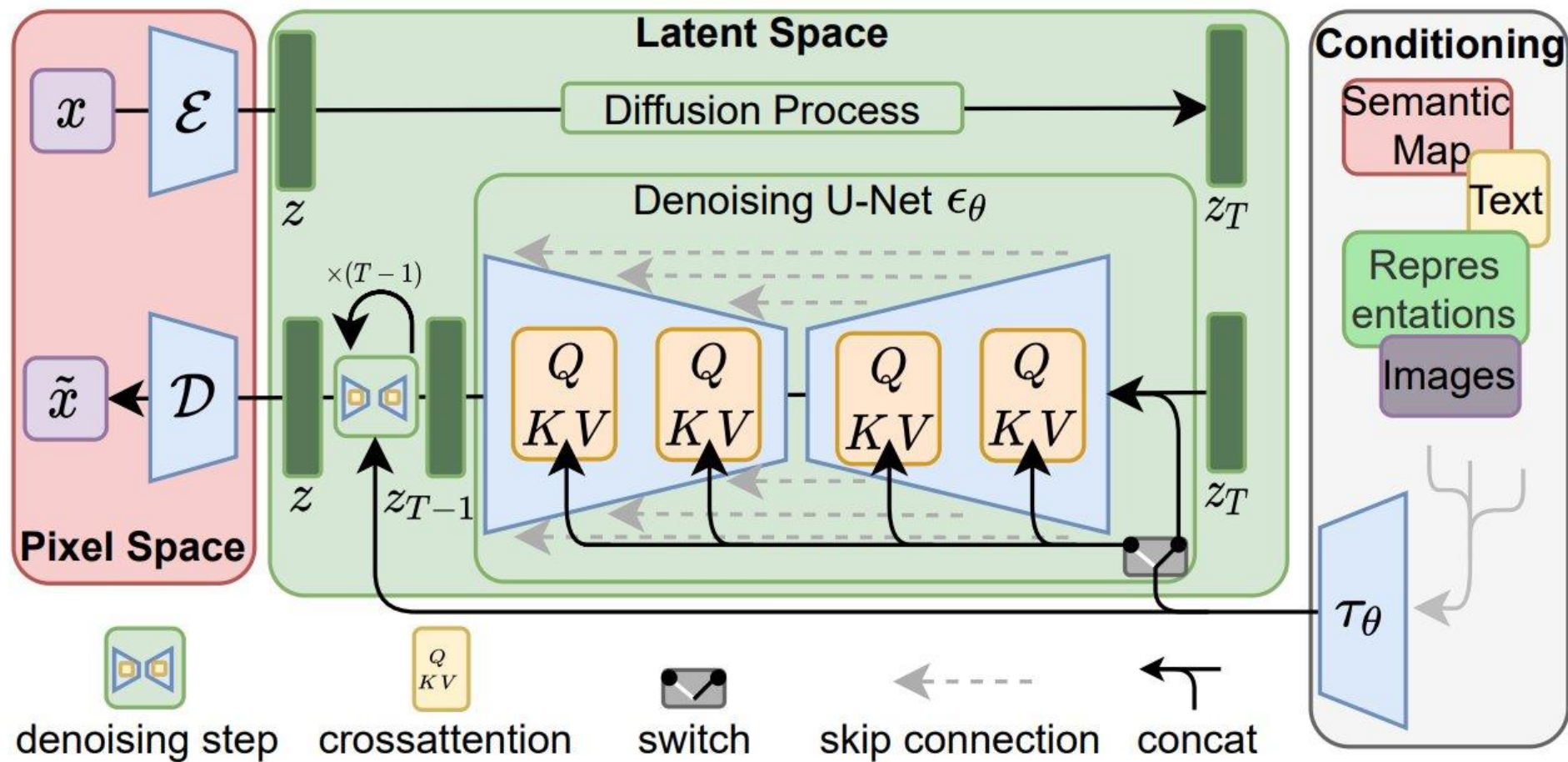
Let the GPUs roar!

Training data & details.

Large Data Training

- SD is trained on ~ 2 Billion image – caption (English) pairs.
 - Scraped from web, filtered by CLIP.
 - <https://laion.ai/blog/laion-5b/>

Recap

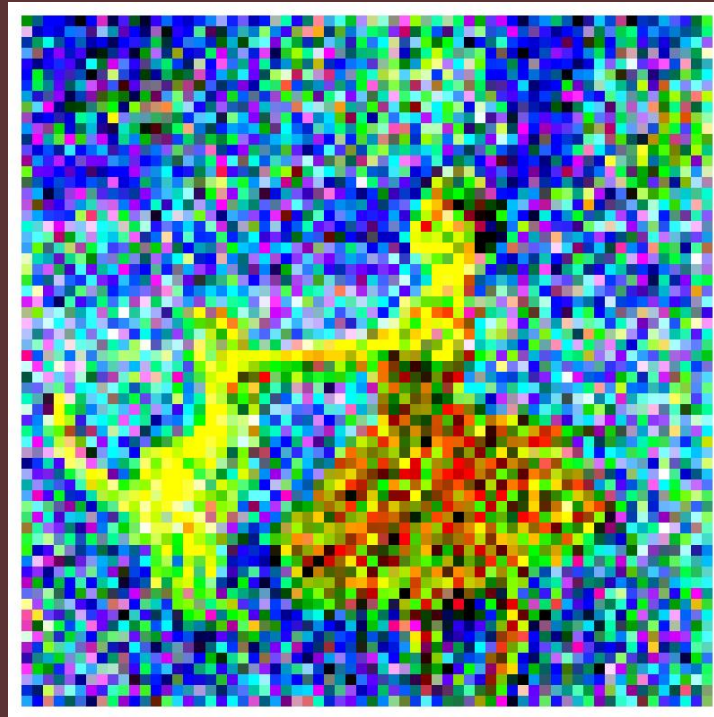


Diffusion Process Visualized



Meaning of latent space

- Latent state contains a “sketch version” of the image.



$z[0:3, :, :]$