Advanced AI: Text to Image Generation

About me









Generative AI | Large Language Models | NLP

United Kingdom · Contact info





■ Live Course



Hands-on Retrieval Augmented Generation (RAG)

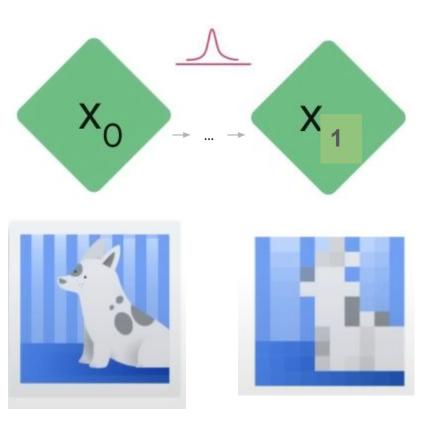
With Jonathan Fernandes

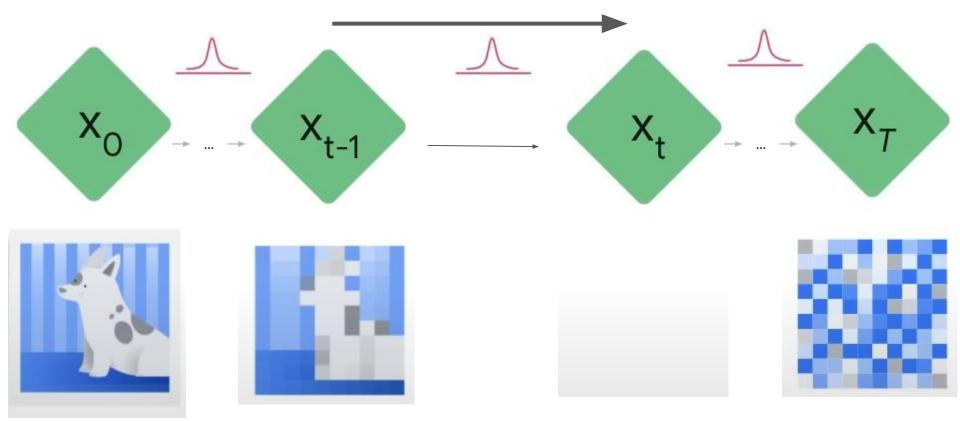
② 3h 0m 🛗 Aug 29 • 5pm-8pm

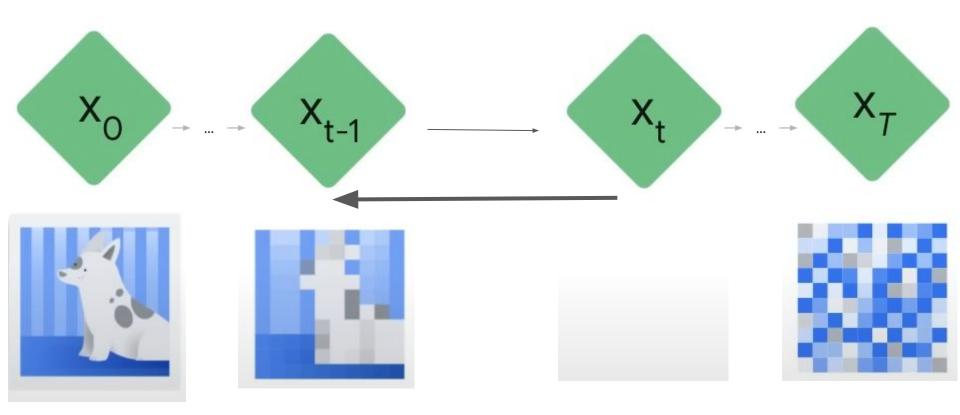
What is diffusion?











Go to notebook

https://playground.com/

https://playground.com/

- (Set the number of images generated to 1 or 2 otherwise, you will run out of credits quickly)
- Shared google doc (image and prompt)
- 3-minute exercise
- Our end goal Use text to generate images

The latest and greatest in text to

image generation

FLUX.1

- Released this month
- Black forest labs
- 3 flavours
 - FLUX.1 dev
 - FLUX.1 schnell
 - FLUX.1 pro
- No NSFW filter
- https://huggingface.co/black-forest-labs/FLUX.1-dev
- https://huggingface.co/spaces/black-forest-labs/FLUX.1-dev

What are some of the problems here?





How would you ensure there is no

nudity/violence/gore in the images that are

generated? [3 minutes]

Can you provide images and associated

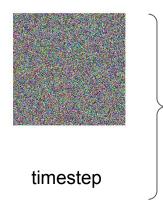
prompts of bias [5 minutes]

Image inpainting:

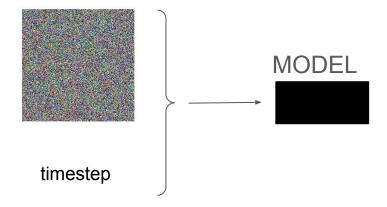
https://huggingface.co/spaces/SkalskiP/FLUX.1-inpaint

[3 minutes]

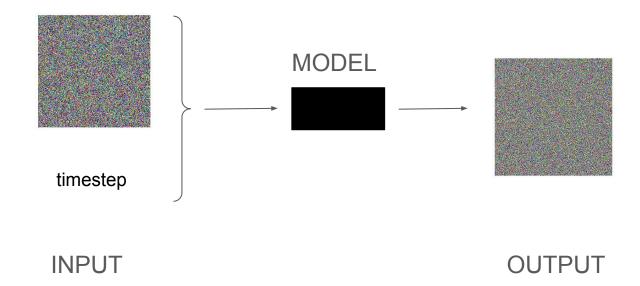
High level overview

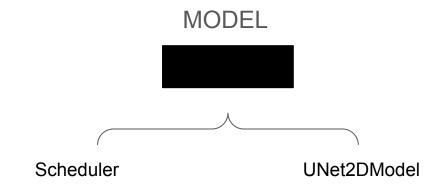


INPUT



INPUT





Go to notebook

Types of Diffusion Models

- Unconditional
- Conditional

Scheduler

$\mathbf{x}_t = \mathbf{x}_{t-1} + \epsilon$

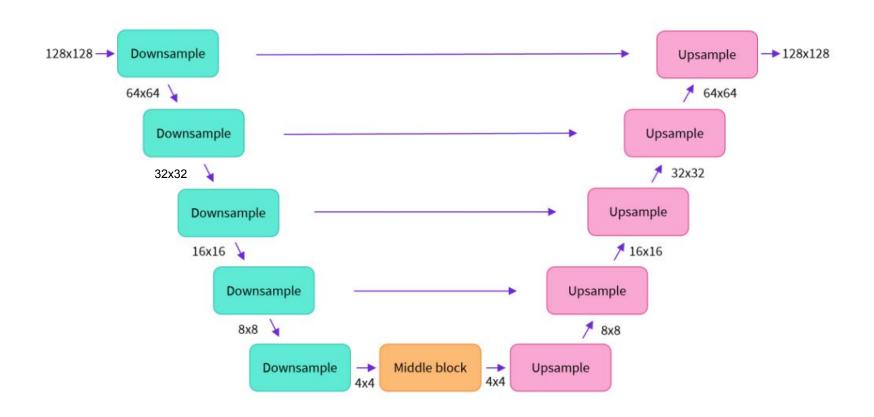
$\mathbf{x}_t = \sqrt{1 - \beta_t} \mathbf{x}_{t-1} + \sqrt{\beta_t} \epsilon$

$$q\left(\mathbf{x}_{t}|\mathbf{x}_{t-1}
ight) = \mathscr{N}\left(\mathbf{x}_{t};\sqrt{1-eta_{t}}\mathbf{x}_{t-1},eta_{t}\mathbf{I}
ight)$$

[Notebook: Scheduler]

U-Net model

U-Net model



[Notebook: U-Net Model]

Train a model

Train a model

Load training image from dataset

Train a model

Load training image from dataset

Add varying noise levels for diverse denoising tasks.

Train a model

Load training image from dataset

Add varying noise levels for diverse denoising tasks.

Input noisy images to the model.

Train a model

Load training image from dataset

Add varying noise levels for diverse denoising tasks.

Input noisy images to the model.

Evaluate model's denoising performance.

Train a model

Load training image from dataset

Add varying noise levels for diverse denoising tasks.

Input noisy images to the model.

Evaluate model's denoising performance.

Update model weights based on evaluation.

Notebook: Train a model

Evaluating Generated Images

Creating Artwork (Generating Images)

- Creating Artwork (Generating Images)
- Art Inspector (Inception Model)

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- Creating Artwork (Generating Images)
- Art Inspector (Inception Model)
- Gallery Walk (Feature Extraction)

- Creating Artwork (Generating Images)
- Art Inspector (Inception Model)
- Gallery Walk (Feature Extraction)
- Comparing Notes (FID score)

Limitations of FID

- Compare distributions, not for a single image
- Not good for low or high-resolution images
- Sensitive to many factors

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- Compare distributions, not for a single image
- Not good for low or high-resolution images
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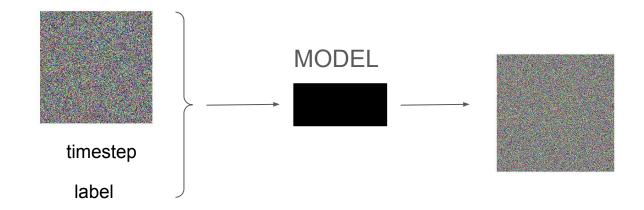
Limitations of FID

- Compare distributions, not for a single image
- Not good for low or high-resolution images
- Sensitive to many factors

Conditioned Diffusion models

https://github.com/zalandoresearch/fashion-mnist

Diffusion models (inference)



INPUT OUTPUT

[Notebook: Conditioned Diffusion Models]

Exercise - CIFAR-10

Solution

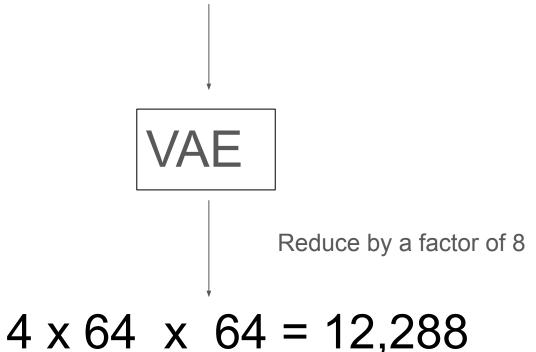
Making improvements - latent diffusion

Latent diffusion

- Scaling
- Images so far have only been 32x32

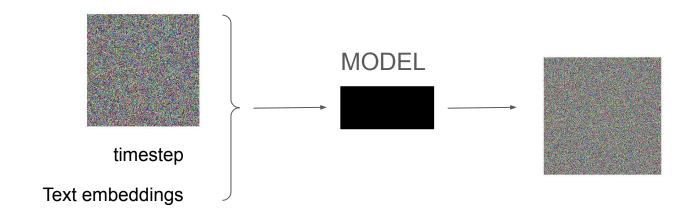
$3 \times 512 \times 512 = 786,432$

$3 \times 512 \times 512 = 786,432$



Text encoder - CLIP

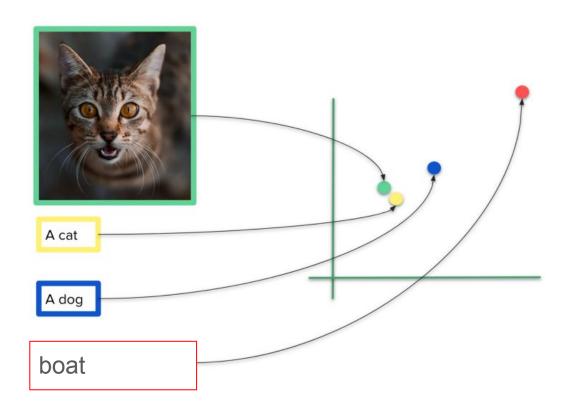
Diffusion models (inference)



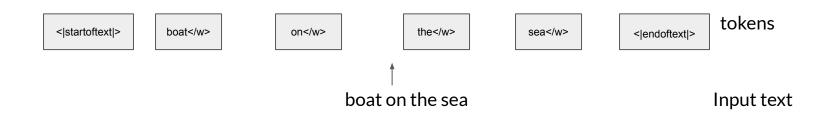
INPUT OUTPUT

https://openai.com/index/clip/

CLIP



boat on the sea

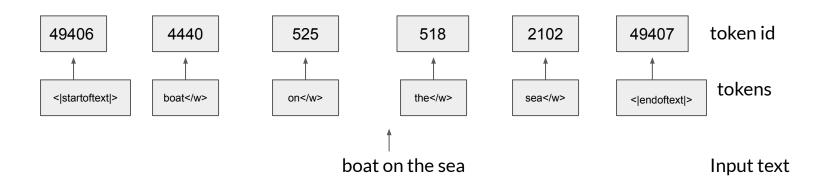


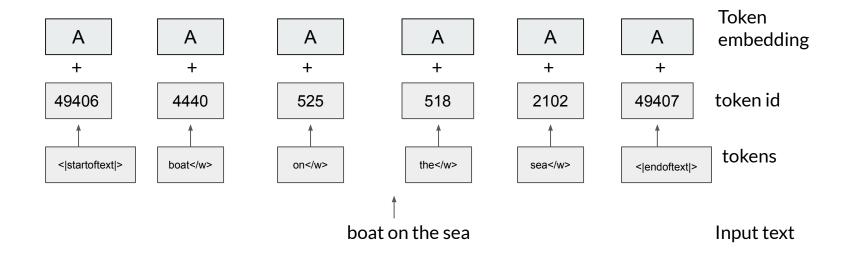
boat on the sea

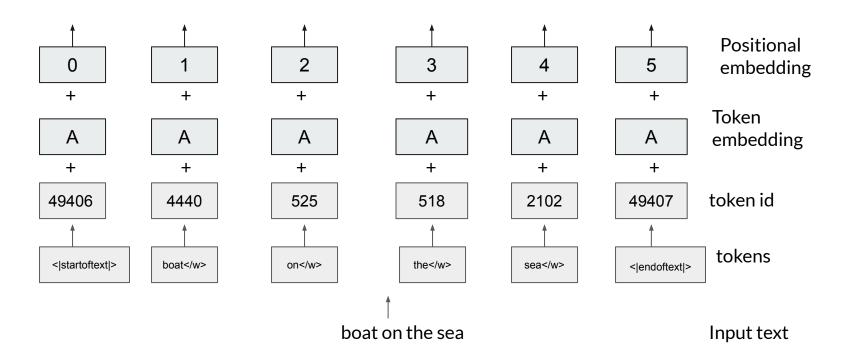
Input text

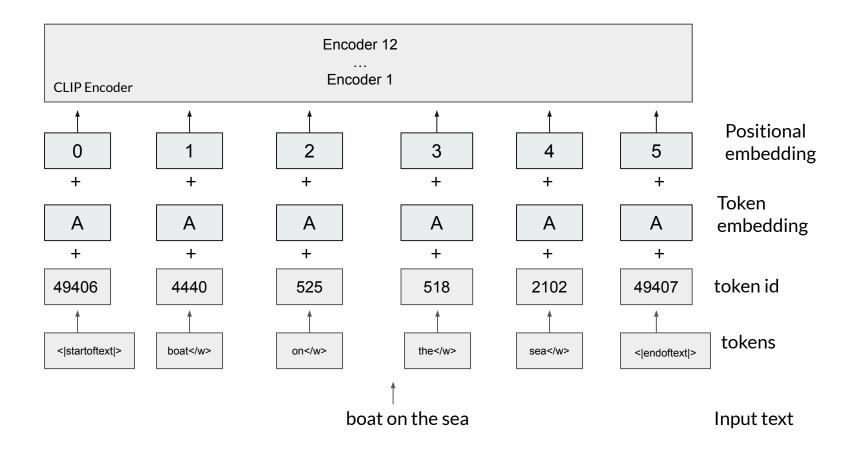
boat on the sea

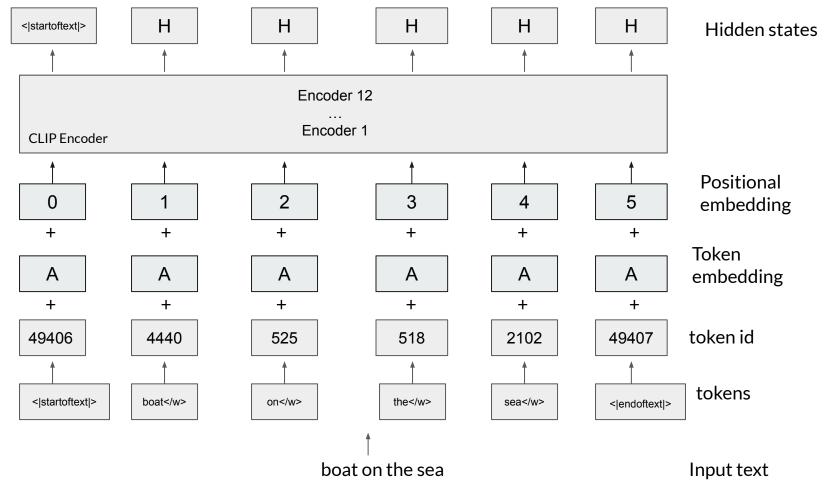
Input text











Text encoder in practice

Notebook: CLIP model

Putting it all together using Stable Diffusion

Bias, Limitations and Controversy

Copyright and Intellectual Property Infringement

- Copyright and Intellectual Property Infringement - Bias and Discrimination

Bias and Discrimination

Copyright and Intellectual Property Infringement

- Deep fakes and miinfornation

Bias and Discrimination

Privacy and consent

- Deep fakes and miinfornation

Copyright and Intellectual Property Infringement

- Copyright and Intellectual Property Infringement
- Bias and Discrimination
- Deep fakes and miinfornation
- Privacy and consent
- Accessibility and control

Next Steps

END

What makes this different to other text to image solutions?

- DALL-E
- DALLE-2
- DALLE-3
- Imagen

Can run on commodity hardware

Model and training details

Model

Text encoder - CLIP ViT-L/14

UNet = 860M parameter model

Autoencoder - downsampling factor of 8.

The model was pretrained on 256x256 images and then finetuned on 512x512 images.

Training time

Hardware Type: A100 PCIe 40GB

Hours used: 150000

Training data

The core dataset was trained on LAION-Aesthetics, a soon to be released subset of LAION 5B.

LAION-Aesthetics was created with a new CLIP-based model that filtered LAION-5B based on how "beautiful" an image was, building on ratings from the alpha testers of Stable Diffusion.

Cost



Jack Clark @jackclarkSF · 28 Aug

Stable Diffusion: \$600k to train.

I'm impressed and somewhat surprised - I figured it'd have cost a bunch more.

Also, Al is going to proliferate and change the world quite quickly if you can train decent generative models with less than \$1m.

...



Emad @EMostaque · 28 Aug

Replying to @KennethCassel

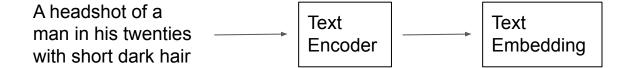
We actually used 256 A100s for this per the model card, 150k hours in total so at market price \$600k

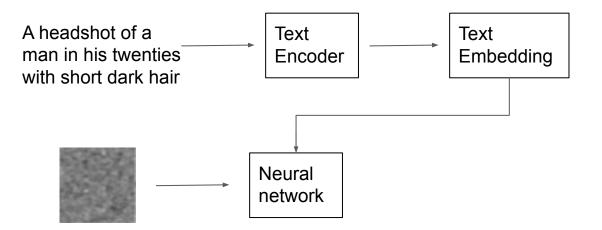
Controversy

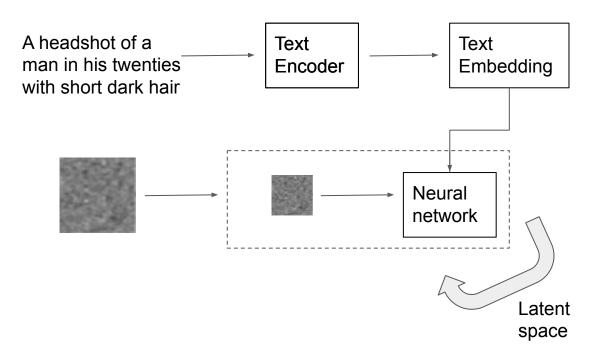
- Image regurgitation
- Copying artist styles
 - Getty Images
 - Shutterstock

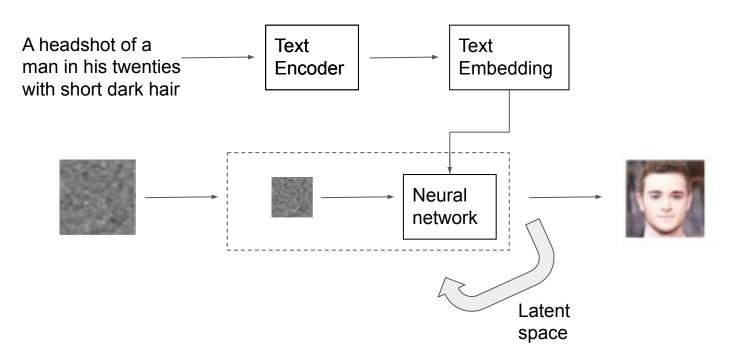
Applications

https://www.reddit.com/r/StableDiffusion/comments/wyduk1/show_rstablediffusion_integrating_sd_in_photoshop/









The key difference between latent and standard diffusion is that latent diffusion model is trained to generate latent (compressed) representations of the images

What 3 components do we need for latent diffusion?

- A text encoder (CLIP's Text Encoder)

What 3 components do we need for latent diffusion?

- A text encoder CLIP's Text Encoder
- Neural network UNet

U-Net



+ Time encoding

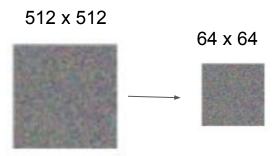
U-Net

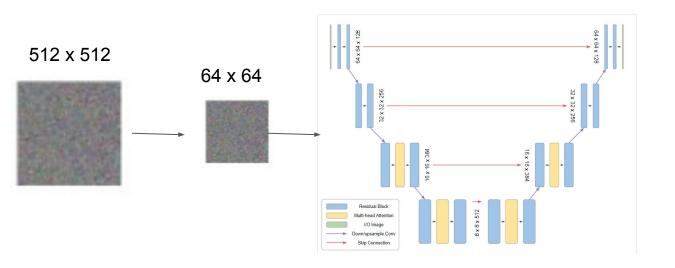


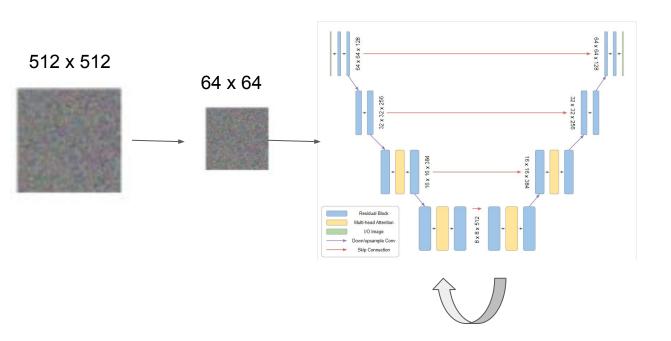
- + Time encoding
- + Text encoding
- Conditioned Image

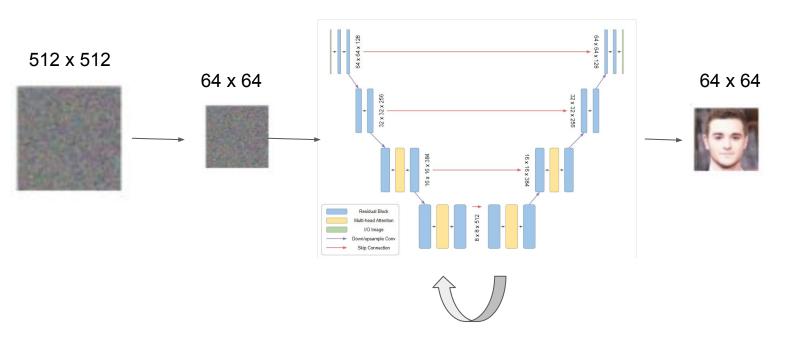
What 3 components do we need for stable diffusion?

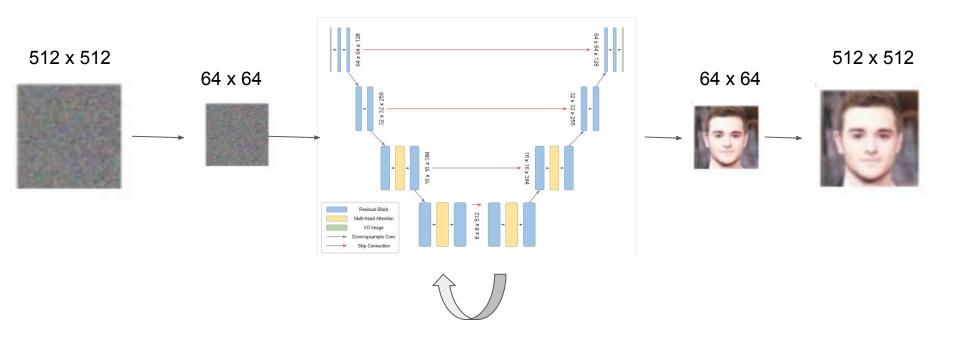
- A text encoder CLIP's Text Encoder
- Neural network UNet
- Autoencoder







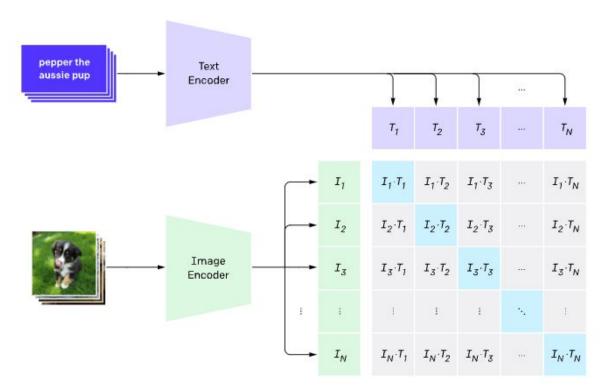




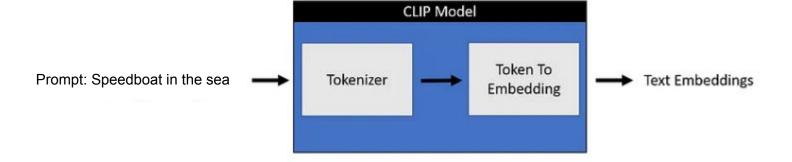
Colab notebook - autoencoder

Text encoder - CLIP

CLIP - Contrastive pre-training



Source: https://openai.com/blog/clip/

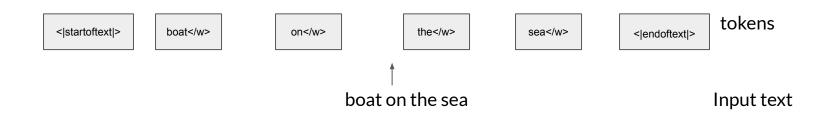


Colab - CLIP Model

Colab - tokenizers

Text encoders

boat on the sea

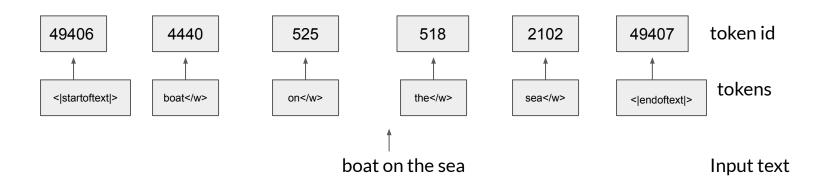


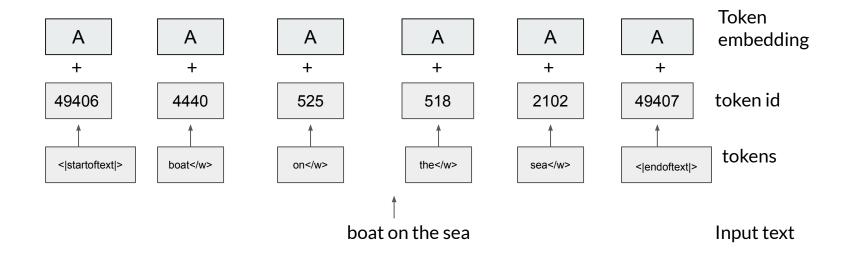
boat on the sea

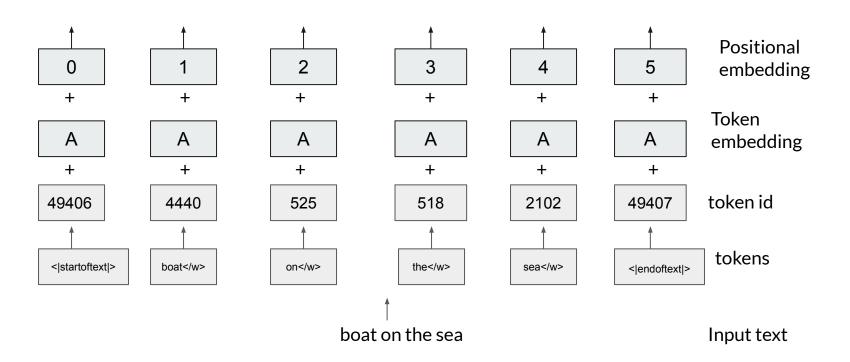
Input text

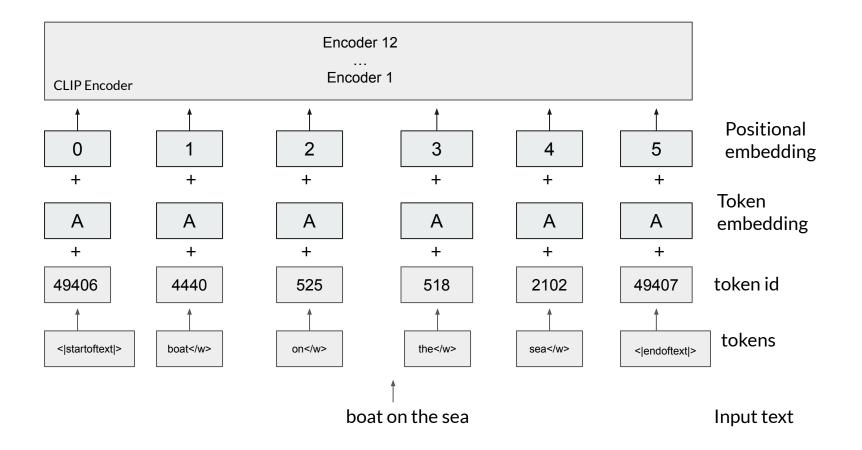
boat on the sea

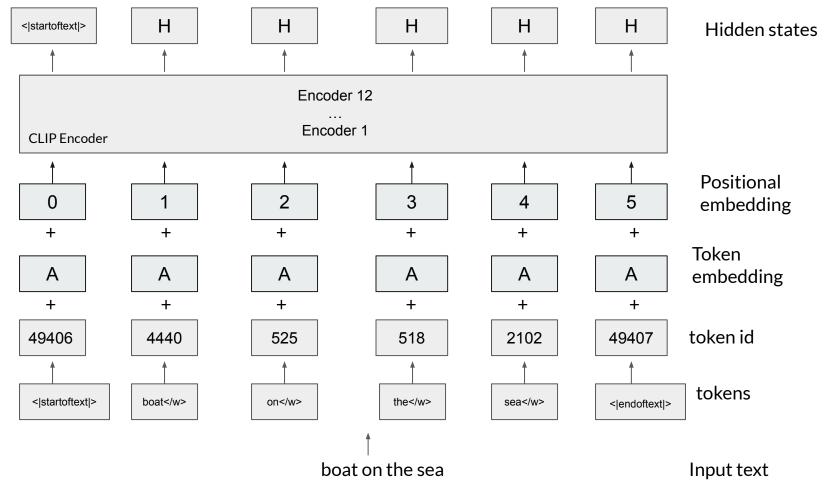
Input text











Colab - text encoders

END