

Controlling Text-to-Image Diffusion Models



BigMAC Workshop @ ICCV'23
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Hugging Face 😊



Disclaimer: This talk is NOT
an exhaustive overview of all
possible methods.

Era of text-to-image diffusion models!



*“A transparent
sculpture of a duck
made out of glass.”*

Imagen



*“panda mad scientist
mixing sparkling
chemicals, digital art.”*

DALL-E 2



*“Astronaut in a jungle,
cold color palette,
muted colors, detailed,
8k”*

SDXL

Diffusion models in a jiffy

What happens when you refine a noise vector to become a realistic image?

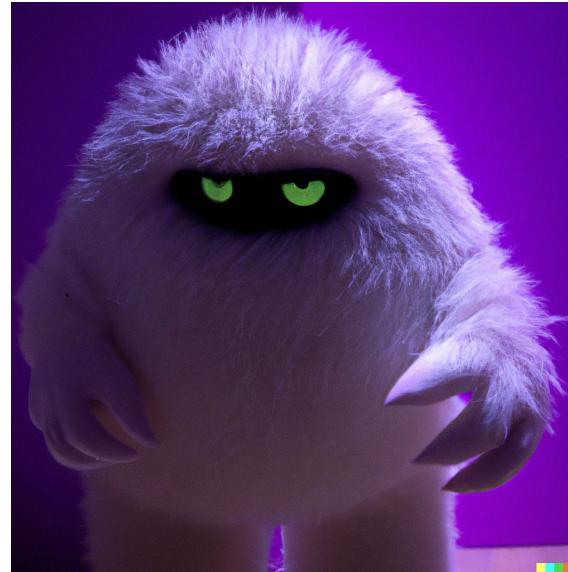
Data

Noise

<https://nvlabs.github.io/denoising-diffusion-gan/>

Diffusion models in a jiffy

When you “condition” the denoising process with text:



DALL-E 2 prompt: “A photo of a white fur monster
standing in a purple room”

Extracting visual connections from textual concepts

Concept discovery in text-to-image diffusion models:

(a)

Concept decomposition with **CONCEPTOR**



Concept: painter

artist sketching abstract
paint Picasso brushes
Pollock Van Gogh ladder
CONCEPTOR →
Monet portrait
staining studio serene
teal Impressionism
calligraphy gilded craftsman

(b) Single-image decomposition with **CONCEPTOR**



Concept: sweet peppers



Concept: beetle

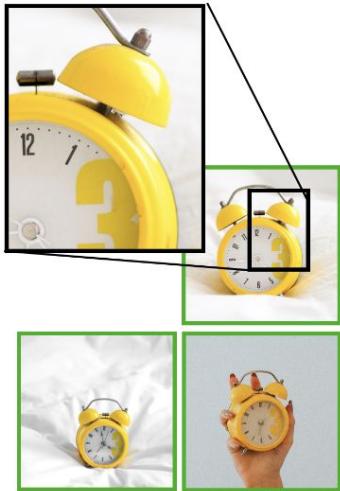


We'll focus on “latent-space” diffusion models throughout this talk. More specifically, the **Stable Diffusion** family.

Limitations and solutions

Part I

Subject-driven generation for personalization

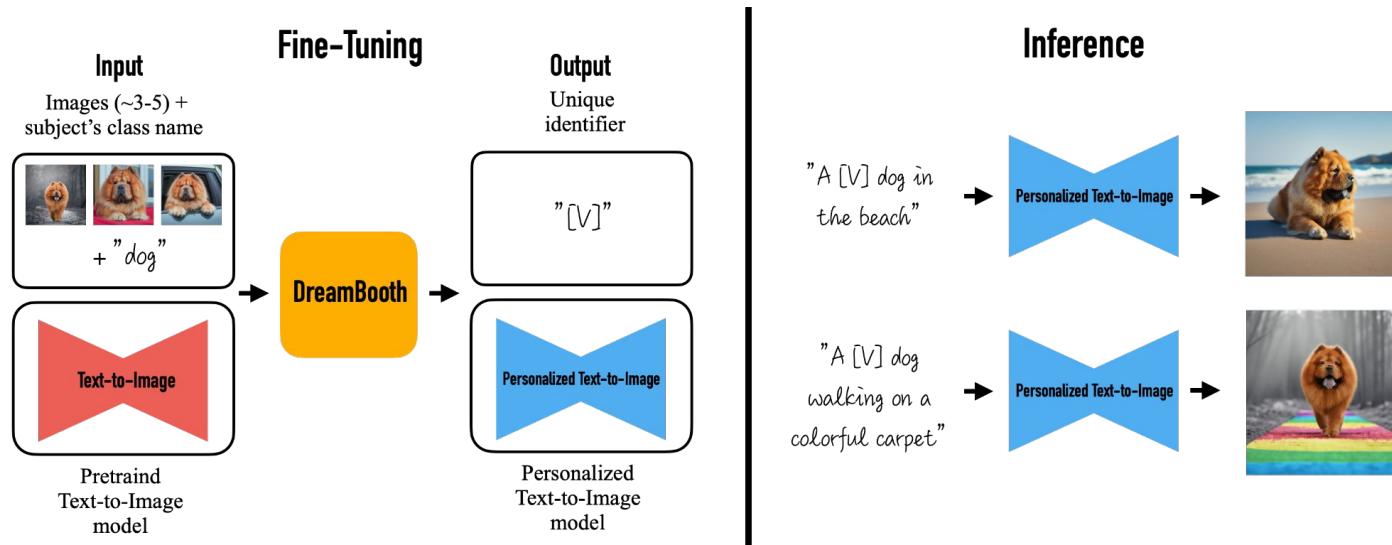


- Render concepts/subjects *new* to the model in interesting contexts.
- Introduce *personalization*.

<https://dreambooth.github.io/>

Subject-driven generation for personalization

Embedding a new subject in the output domain of the (pre-trained) model:
DreamBooth!



<https://dreambooth.github.io/>

DreamBooth; Ruiz et al., 2022.

Subject-driven generation for personalization

Without the loss of generality, let:

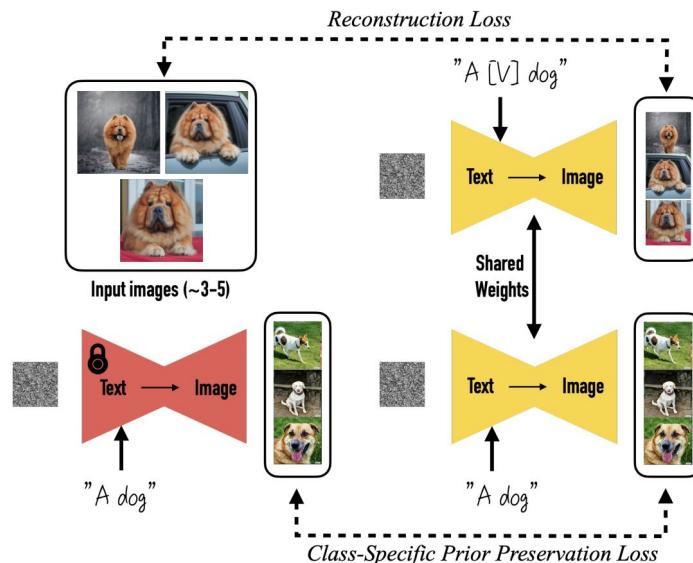
- \mathbf{x} : original image
- $\boldsymbol{\epsilon}$: noise; $\boldsymbol{\epsilon} \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$
- t : diffusion process time; $t \sim \mathcal{U}([0, 1])$
- a_t, σ_t, w_t : terms controlling noise schedule and sample quality
- \mathbf{c} : conditioning vector (prompt embeddings, for example)
- $\hat{\mathbf{x}}_\theta$: diffusion model to be learned

Training $\mathbb{E}_{\mathbf{x}, \mathbf{c}, \boldsymbol{\epsilon}, t} \left[w_t \|\hat{\mathbf{x}}_\theta(a_t \mathbf{x} + \sigma_t \boldsymbol{\epsilon}, \mathbf{c}) - \mathbf{x}\|_2^2 \right]$

Inference $\mathbf{x}_{\text{gen}} = \hat{\mathbf{x}}_\theta(\boldsymbol{\epsilon}, \mathbf{c})$

Subject-driven generation for personalization

Prior-preservation loss to preserve the class-specific semantic prior:

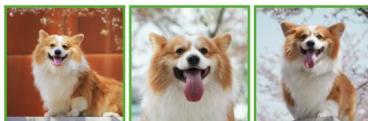
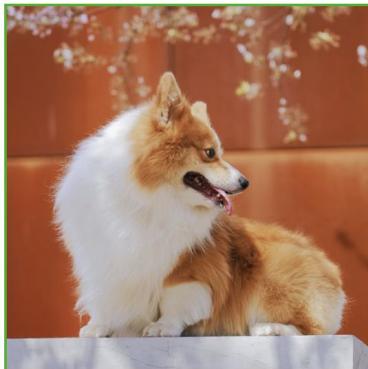


$$\mathbb{E}_{\mathbf{x}, \mathbf{c}, \epsilon, \epsilon', t} [w_t \|\hat{\mathbf{x}}_\theta(\alpha_t \mathbf{x} + \sigma_t \epsilon, \mathbf{c}) - \mathbf{x}\|_2^2 + \lambda w_{t'} \|\hat{\mathbf{x}}_\theta(\alpha_{t'} \mathbf{x}_{\text{pr}} + \sigma_{t'} \epsilon', \mathbf{c}_{\text{pr}}) - \mathbf{x}_{\text{pr}}\|_2^2]$$

One framework, multiple use cases

General subject-driven generation

Input images



A [V] dog in the
Versailles hall of mirrors



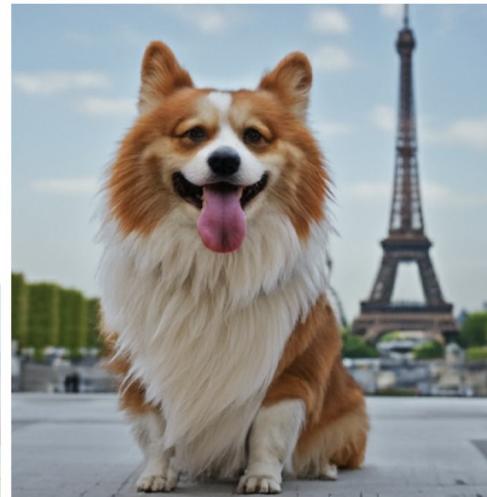
A [V] dog in the
gardens of Versailles



A [V] dog in Coachella



A [V] dog in
mountain Fuji

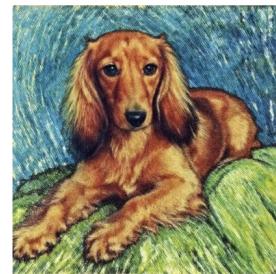


A [V] dog with Eiffel Tower in
the background

One framework, multiple use cases

Art rendition

Input images



Vincent Van Gogh



Michelangelo



Rembrandt



Johannes Vermeer



Pierre-Auguste Renoir



Leonardo da Vinci

One framework, multiple use cases

Property modification

Hybrids (“A cross of a [V] dog and a [target species]”)



Input



Bear



Panda



Koala



Lion



Hippo

Pushing the extremes with DreamBooth

Order by random likes

SDXL LoRA Gallery

Click on a LoRA in the gallery to select it

Type a prompt after selecting a LoRA

Run

Generated Image

Advanced options

pe-lofi-hiphop-lofi-girl-c...

Voxel XL

Lego BrickHeadz

pe-funko-pop-diffusion-...

1987-action-figure-play...

CAG Coinmaker

PixelArtRedmond

Toy.Redmond

toy-face



Further reads

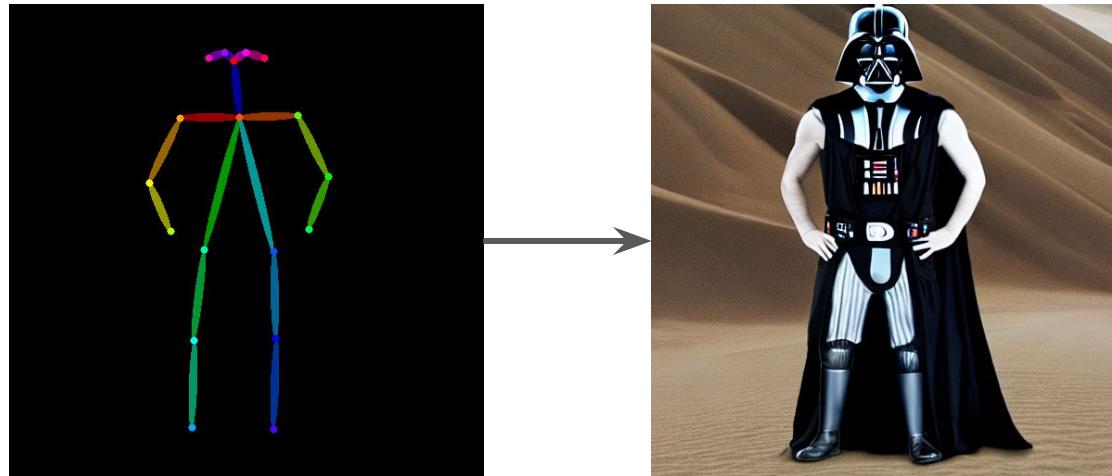
- **BLIP-Diffusion**; Li et al., 2023 (zero-shot subject-driven generation).
- **Custom Diffusion**; Kumari et al., 2022.
- **Pivotal Tuning**; Roich et al., 2021 (in SD context it's Textual Inversion + DreamBooth).

Limitations and solutions

Part II

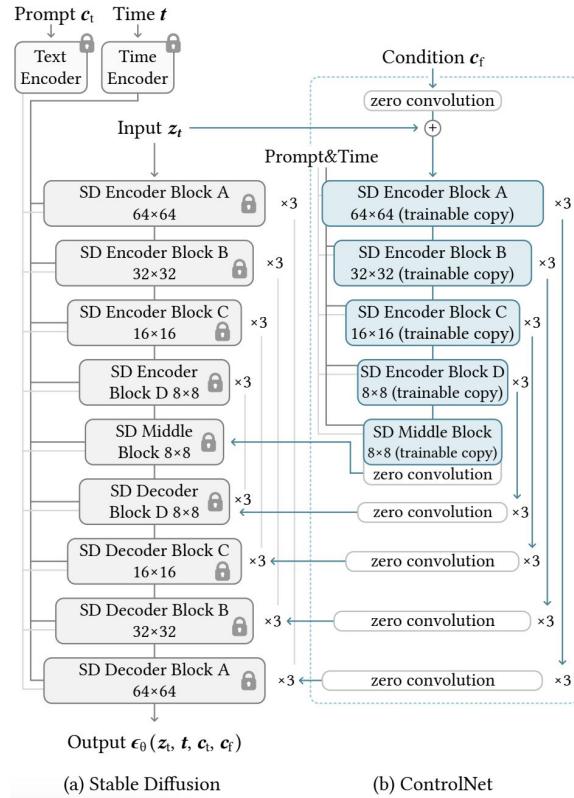
Going beyond text conditioning

What if we wanted to condition the generation process on a pose image along with language supervision?



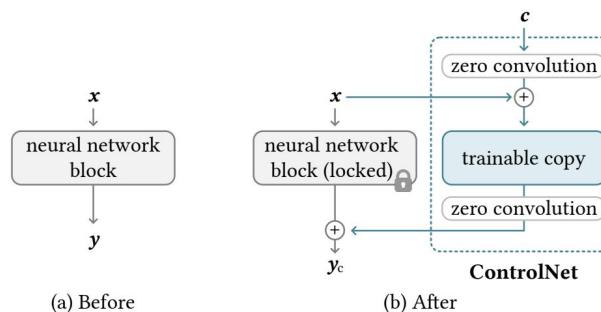
"Darth Vader dancing in a desert"

Going beyond text conditioning - ControlNets

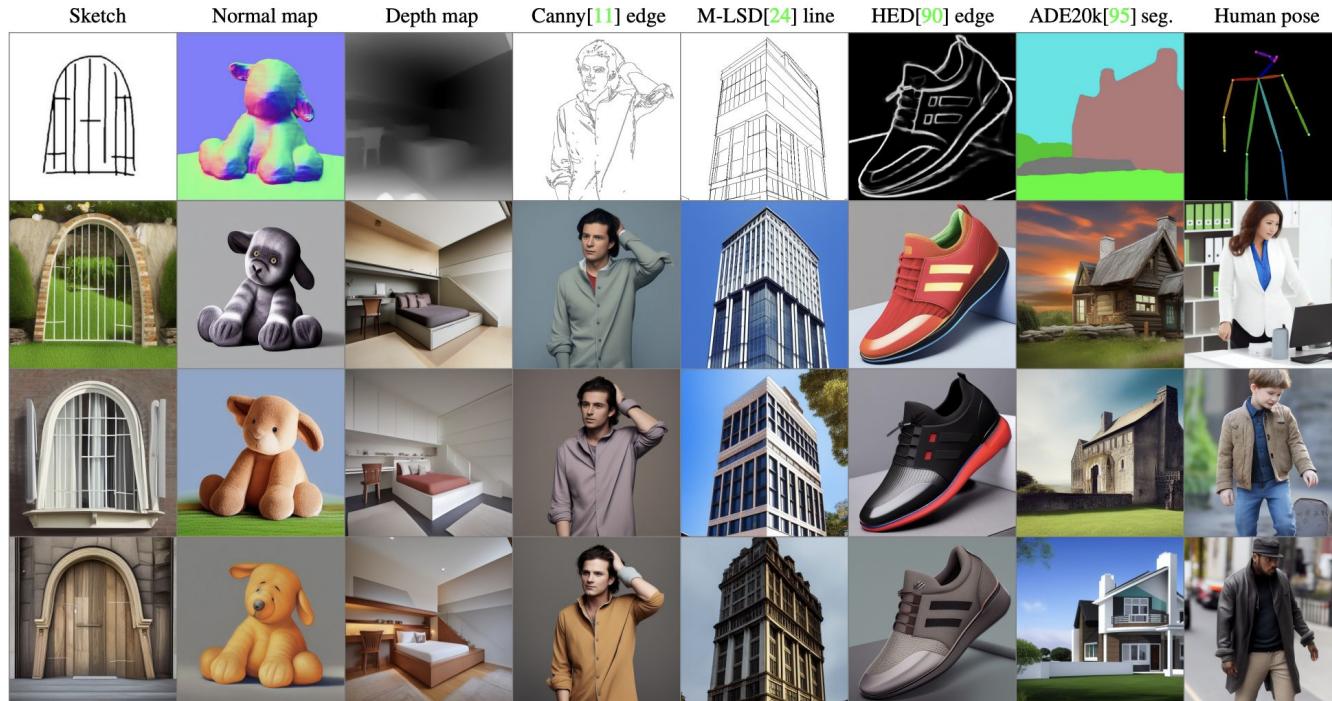


$$\mathbb{E}_{\mathbf{x}, \mathbf{c}_t, \mathbf{c}_f, \boldsymbol{\epsilon}, t} \left[w_t \|\hat{\mathbf{x}}_\theta(\alpha_t \mathbf{x} + \sigma_t \boldsymbol{\epsilon}, \mathbf{c}_t, \mathbf{c}_f) - \mathbf{x}\|_2^2 \right]$$

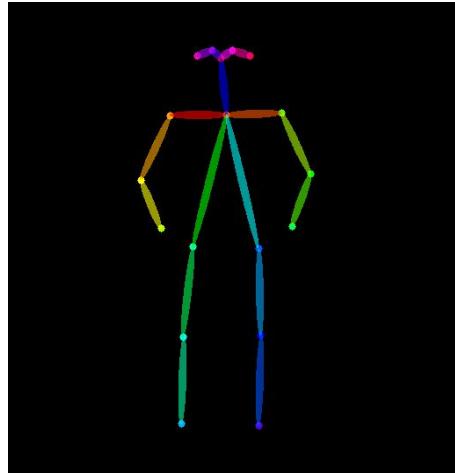
Image-space conditioning vector for ControlNet



ControlNets - a powerful framework to inject additional control



Or shall I say **controls**?



Canny map

Pose

Final Image

*"a giant standing in a fantasy landscape,
best quality"*

Further reads

- **T2I-Adapters**; Mou et al., 2023.
- **IP-Adapters**; Ye et al., 2023.
- **InstructPix2Pix**; Brooks et al., 2022.

Limitations and solutions

Part III

Catastrophic neglect & incorrect attribute binding

“A yellow **bowl** and a blue **cat**”



Neglects one or more objects in the generation.

“A yellow **bow** and a brown **bench**”



Fails to properly bind attributes to objects.



y tho?

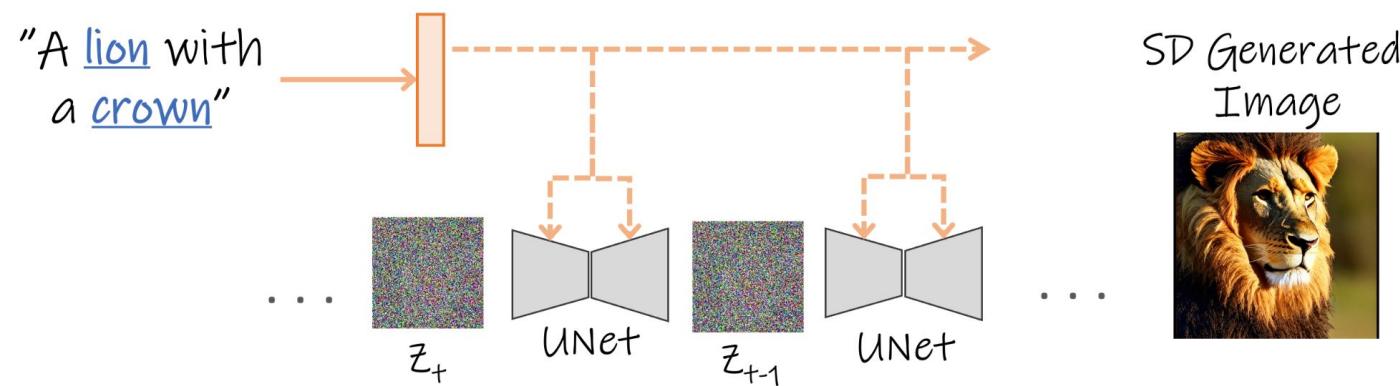


Going to steal a couple of slides
from Hila Chefer here.

Why Does the Model Fail?

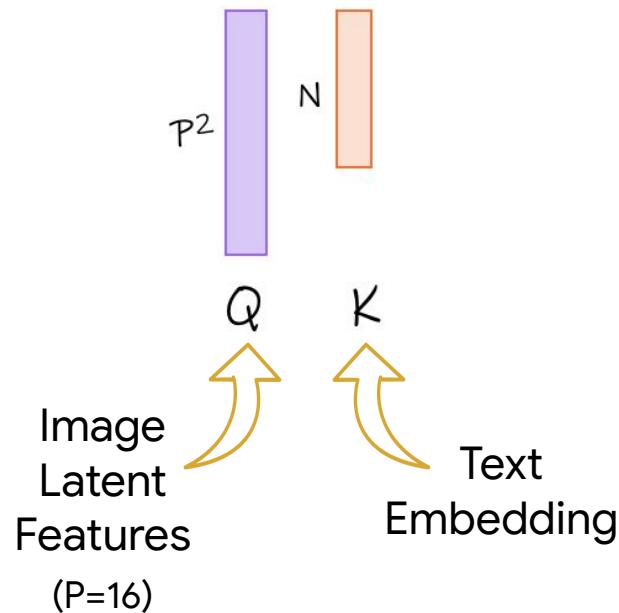
DDPM process:

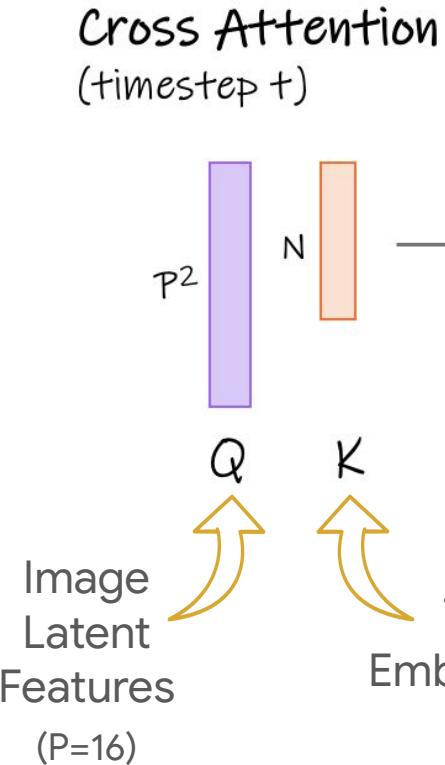
Given an input text prompt, the DDPM gradually denoises a pure noise latent to obtain the output image.



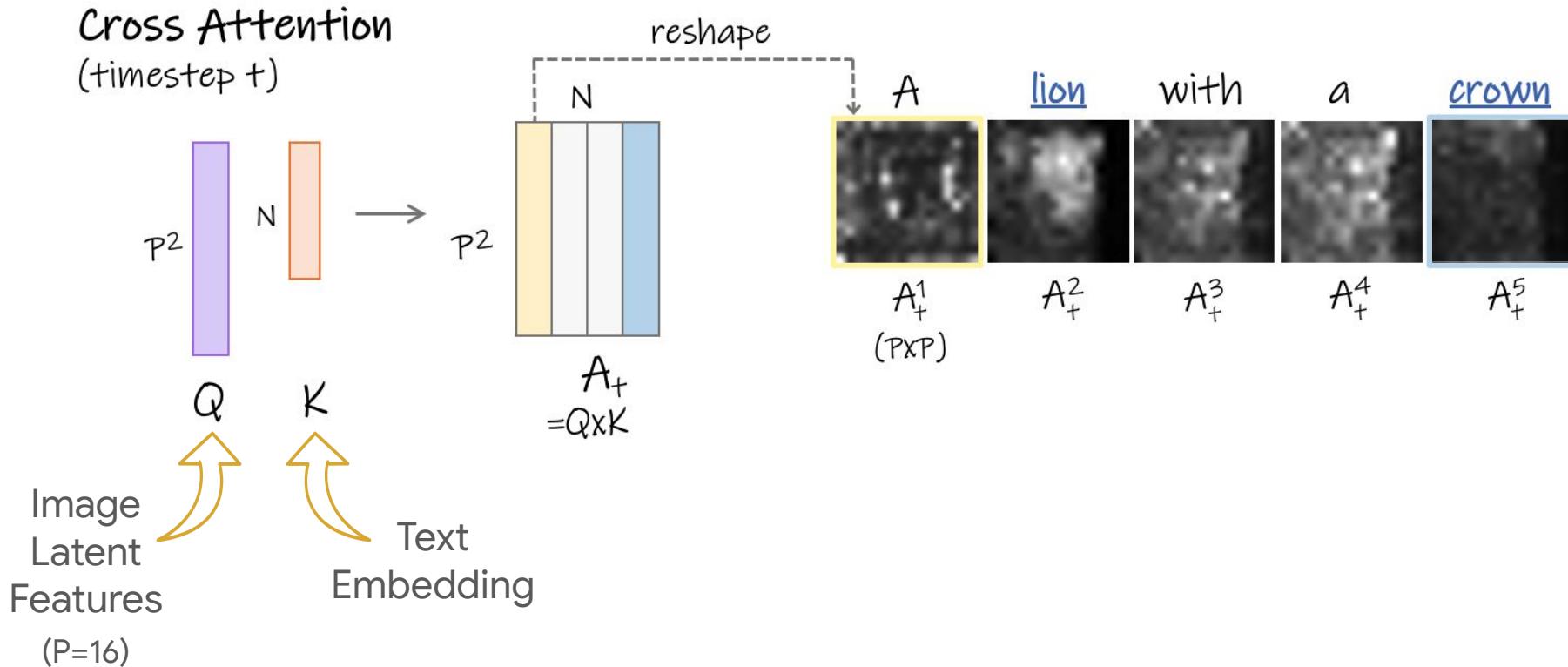
Cross Attention
(timestep +)

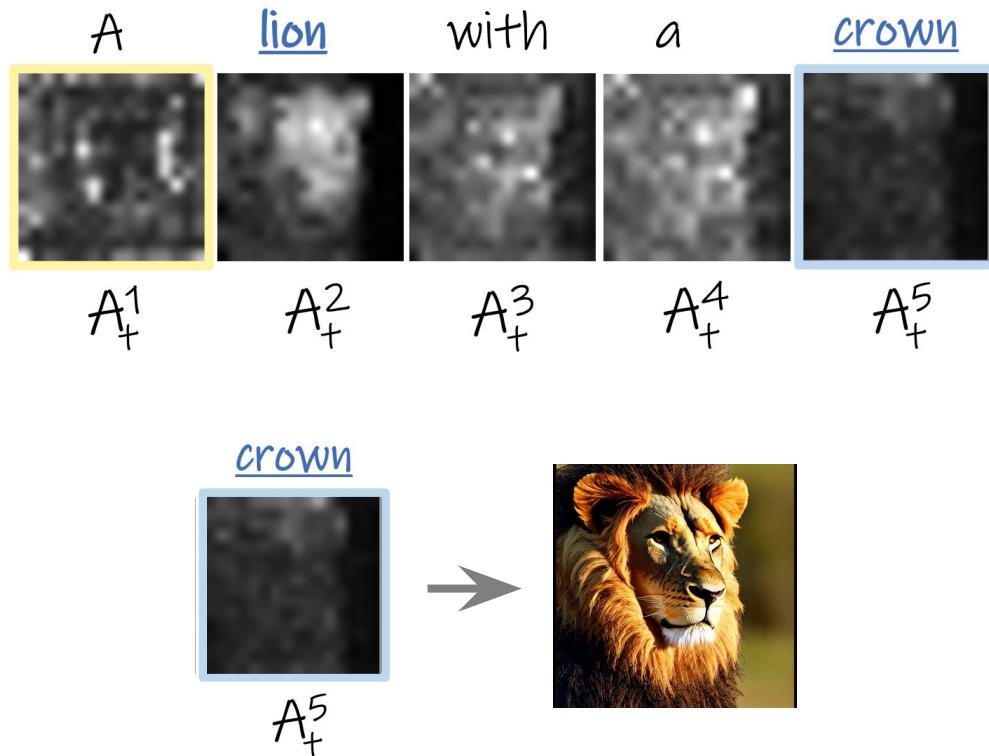
Cross Attention (timestep +)





$A_+[i,n]$ = presence of the
token n in patch i





Problem: crown gets low attention values for all patches

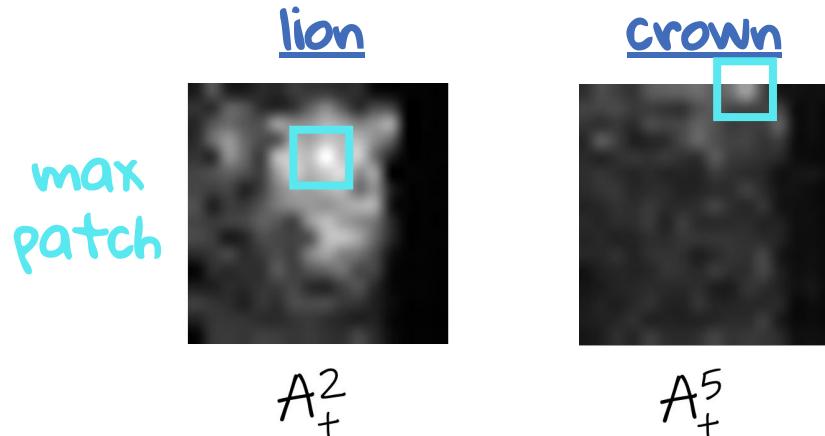
Generative semantic nursing

We want to:

- Encourage the model to better consider the semantic information passed from the input text prompt.
- Ensure all tokens are attended to by some image patch meaningfully.

How can we fix this?

 **Intuition:** a generated subject should have an image patch that significantly attends to the subject's token.



$$L_2 = 1 - \max A_+^2 \quad L_5 = 1 - \max A_+^5$$

Loss: $L = \max(L_2, L_5)$

Update: $\bar{z}_+ = z_+ - \alpha \nabla_{z_+} L$

How close are we to having a strong patch?

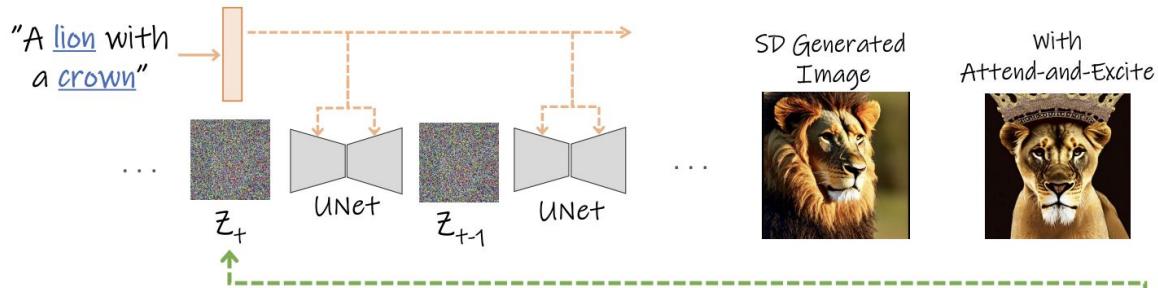


Idea: strengthen the activation of the *most neglected* token

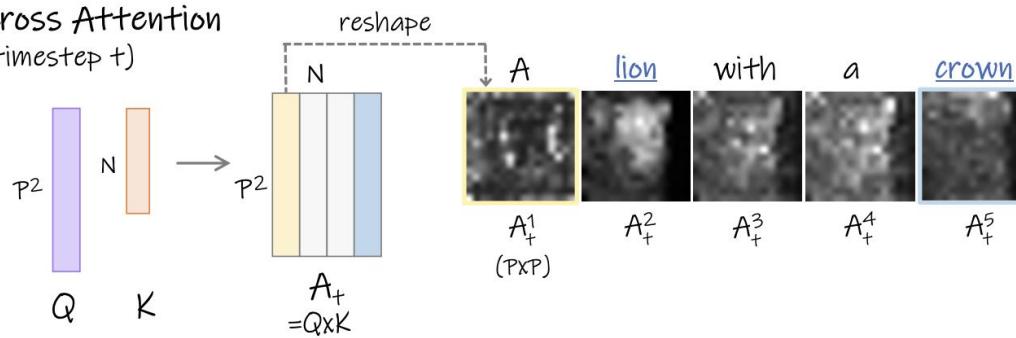
Putting It All Together

Attend to and
Excite all subject
tokens!

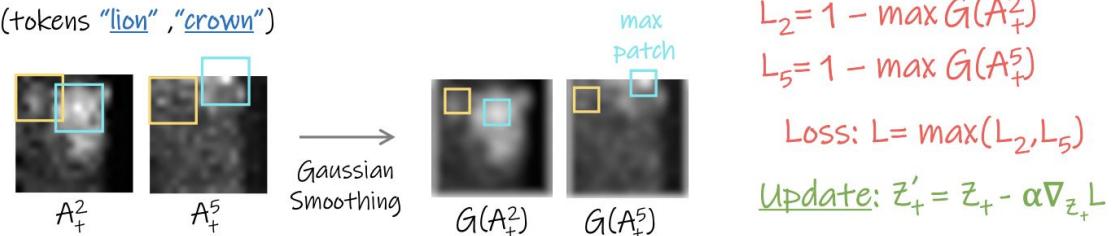
DDPM Process



Cross Attention
(timestep t)



Loss Computation
(tokens "lion", "crown")



Results

"A playful kitten chasing a butterfly in a wildflower meadow"



Stable Diffusion



Attend-and-Excite

Results

"A grizzly bear catching a salmon in a crystal clear river surrounded by a forest"



Stable Diffusion



Attend-and-Excite

Notable mentions

Controlling semantic attributes (training-free):

- Semantic Guidance; Brack et al., 2023.
- LEDITS; Tsaban et al., 2023.

Controlling using “rich-text” (training-free):

- Expressive Text-to-Image Generation with Rich Text; Ge et al., 2023.

Improving discriminative performance:

- Synthetic Data from Diffusion Models Improves ImageNet Classification; Azizi et al., 2023.



IF prompt: A cute panda standing amidst a mountain and holding a placard saying “Thank you!”

