# DragDiffusion: Harnessing Diffusion Models for Interactive Point-based Image Editing

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arXiv(CVPR), 2023(2024), 60 citation

# Interactive point-based image editing

- Move several handle points to target points
- Requirements of task
  - Flexibility: adjust spatial attributes (e.g., pose, position, expression, shape, etc)
  - Precision: control spatial attributes with high precision
  - Generality: apply various categories

Handle point
Target point



# Existing approach

- DragGAN (2023)
  - Due to capacity of GAN model, generality is not satisfied in DragGAN



# Existing approach

- Large-scale text-to-image diffusion models
  - Have a strong capabilities
  - Most diffusion-based editing models use text embeddings
  - It cannot achieve precise spatial control



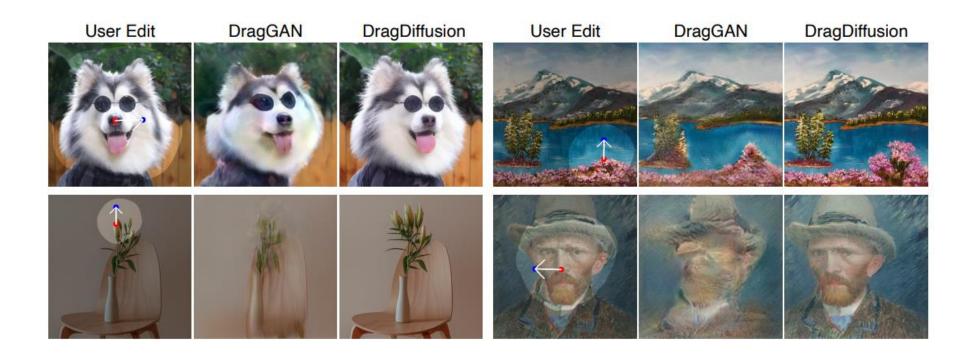


Prompt-to-Prompt

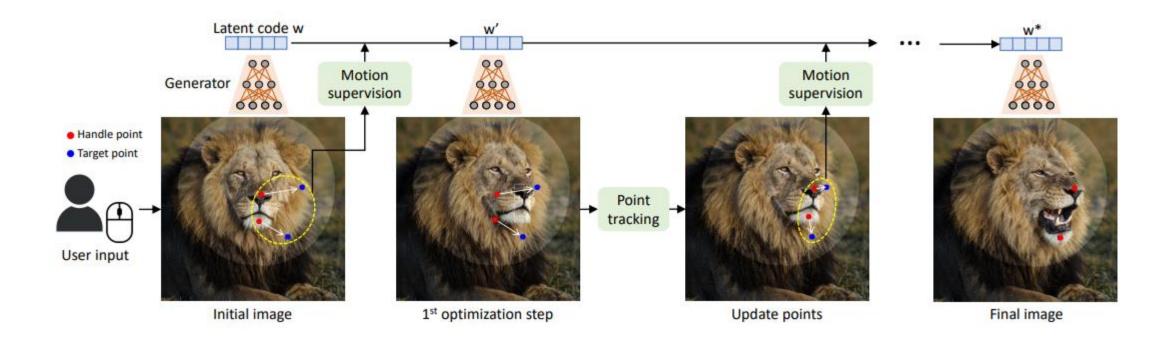
InstructPix2Pix

# DragDiffusion

- Use diffusion models instead of GAN
- Additionally, introduce two techniques to preserve identity of original image
  - Identity-preserving fine-tuning
  - Reference-latent-control



- DragGAN: image manipulation via optimizing latent code
  - StyleGAN2, feature map of 6<sup>th</sup> block
  - Motion supervision loss: move handle point to target point
  - Point tracking loss: update previous handle point to current handle point



- DragGAN: image manipulation via optimizing latent code
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$$\mathcal{L} = \sum_{i=0}^{n} \sum_{q_i \in \Omega_1(p_i, r_1)} \|F(q_i) - F(q_i + d_i)\|_1 + \lambda \|(F - F_0) \cdot (1 - M)\|_1,$$

Motion supervision

$$\mathbf{p}_i \coloneqq \underset{\mathbf{q}_i \in \Omega_2(\mathbf{p}_i, r_2)}{\operatorname{arg \, min}} \|\mathbf{F}'(\mathbf{q}_i) - f_i\|_1.$$

Point tracking

*F*: *feature map* 

 $F_0$ : initial feature map

i: number of points

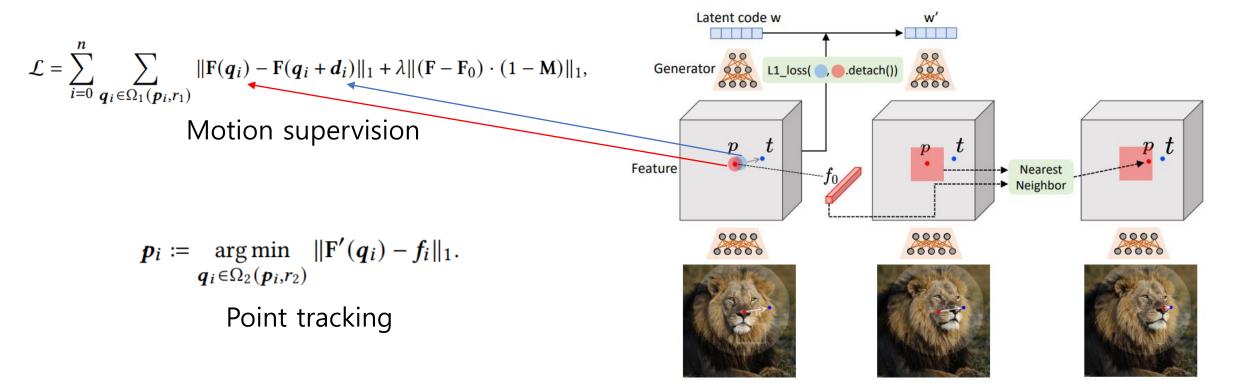
 $d_i$ : unit vector towards target points

 $q_i$ : small patch around  $p_i$ 

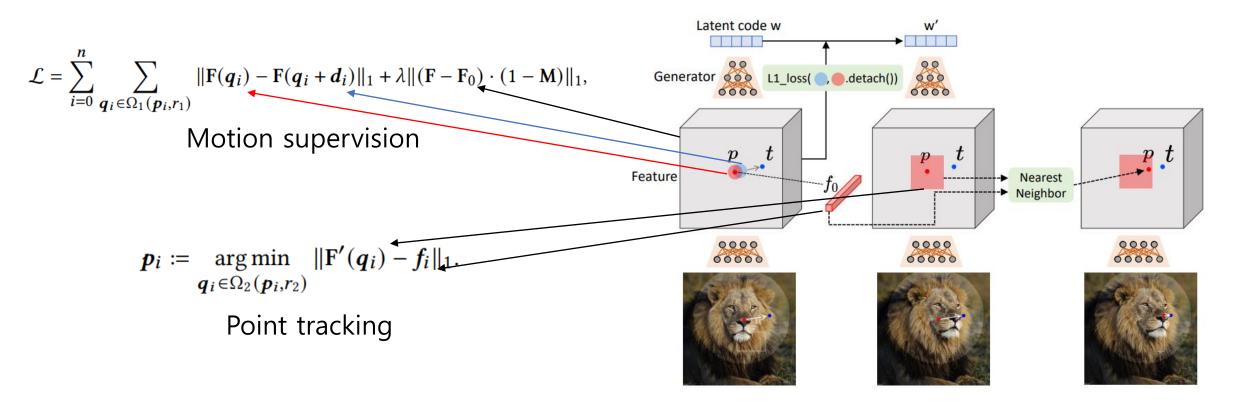
 $f_i$ : feature of initial handle point

 $F'(q_i)$ : updated feature map at  $q_i$ 

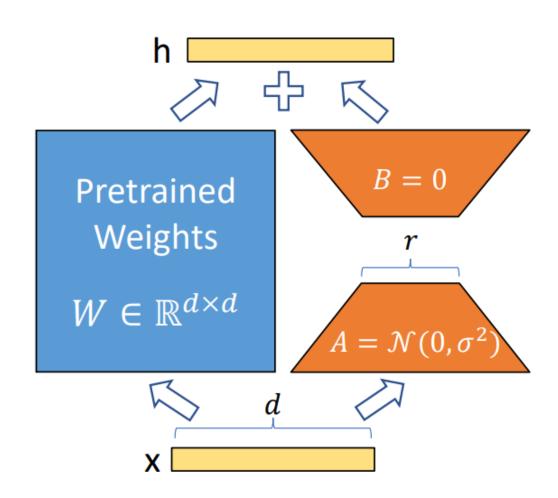
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- LoRA
  - Train the model with additional parameters



- Diffusion models
  - Forward process: add noise / Reverse process: remove noise
  - Given data  $X_0$ , add noise  $\epsilon \sim N(0, I)$  iteratively (forward process)
  - The model trained the forward process in reverse (reverse process)
  - Therefore, the model predicts the noise at a specific time step

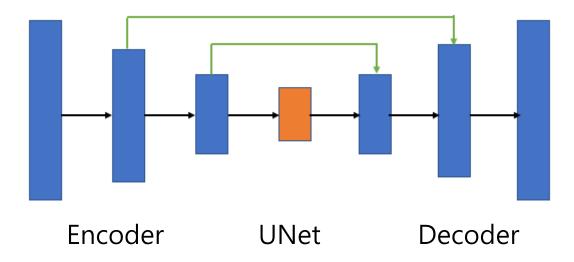
$$x_t = \sqrt{\overline{\alpha_t}} x_0 + \sqrt{(1 - \overline{\alpha_t})} \epsilon$$

Forward process (add noise)

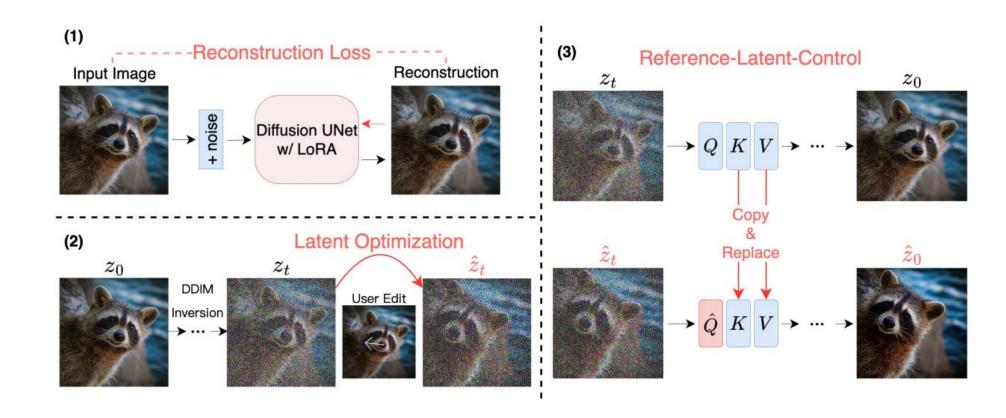


Reverse process (remove noise)

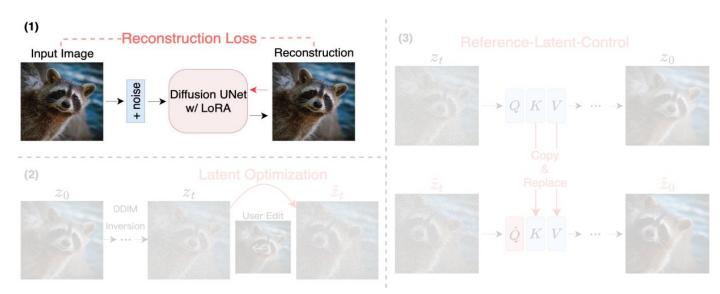
- Diffusion models
  - Forward process: add noise / Reverse process: remove noise
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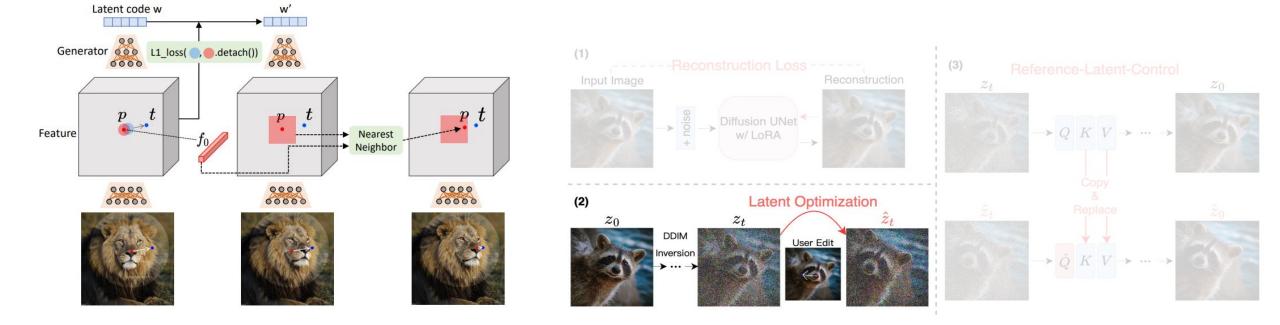
- DragDiffusion
  - Use diffusion model
  - It consists of 3 step



- 1. Identity-preserving fine-tuning
  - To encode input's feature in diffusion model, train model using LoRA
  - Only 80 steps
    - Subject-driven image generation: require 1000 steps (DreamBooth, textual inversion)
    - Sampling time: GAN < Diffusion + fine-tuning
    - A100GPU: 25 seconds



- 2. Diffusion latent optimization (motion supervision + point tracking)
  - Move handle point to target point + update handle point
  - Same with DragGAN



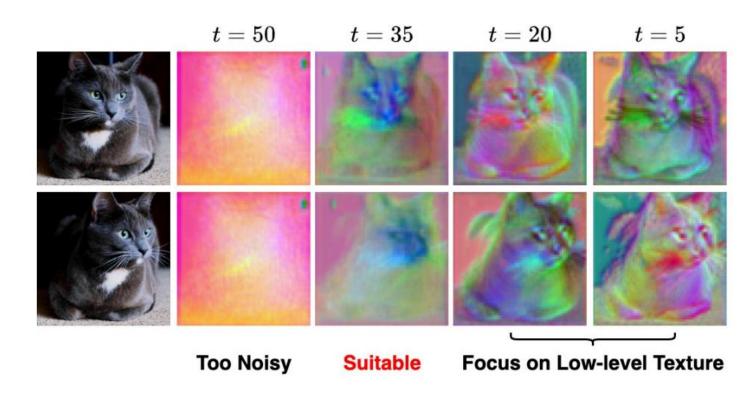
2. Diffusion latent optimization (motion supervision + point tracking)

- GAN: generate an image at once
- Diffusion: generate an image with iterative denoising

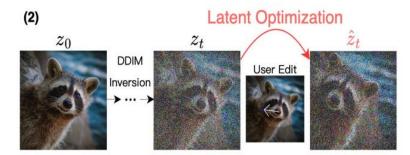


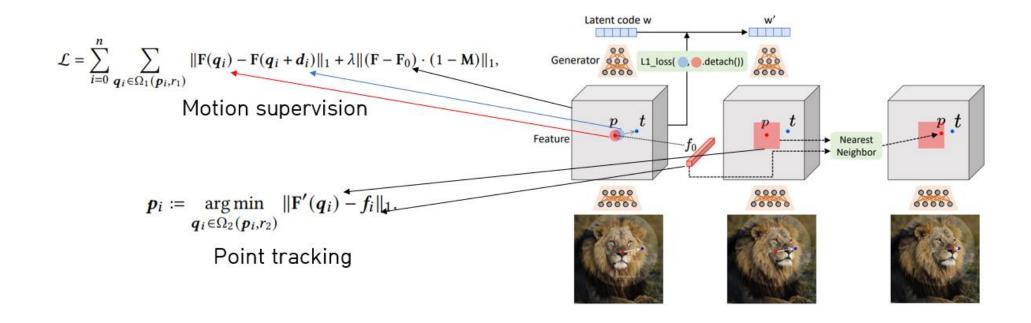
All of time steps? Certain time step?

- 2. Diffusion latent optimization (motion supervision + point tracking)
  - Given two frames, visualize feature map over time using PCA
  - At t = 35, it has sufficient semantic and geometric information (shape, pose, etc)
  - Conduct optimization at certain time step (t = 35)



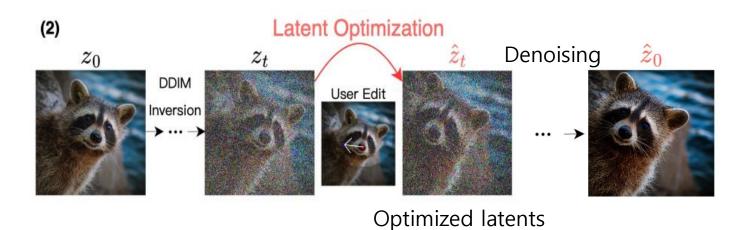
- 2. Diffusion latent optimization (motion supervision + point tracking)
  - Optimization step: 80





#### 3. Reference-latent-control

- After optimization, denoise optimized latents to generate final editing results
- Occurs defects: shift, degrade quality
- Assume that this issue arises due to the absence of proper guidance from the original image during denoising process

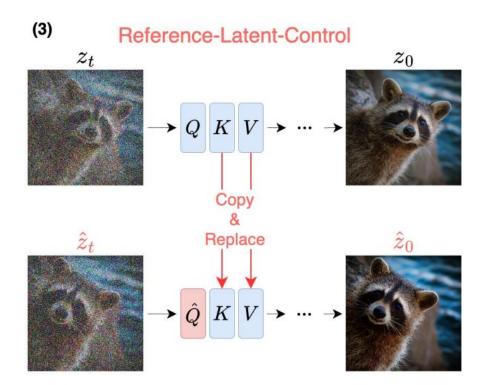


#### 3. Reference-latent-control

 In self-attention module, replace key and value of optimized latents with key and value of original latents

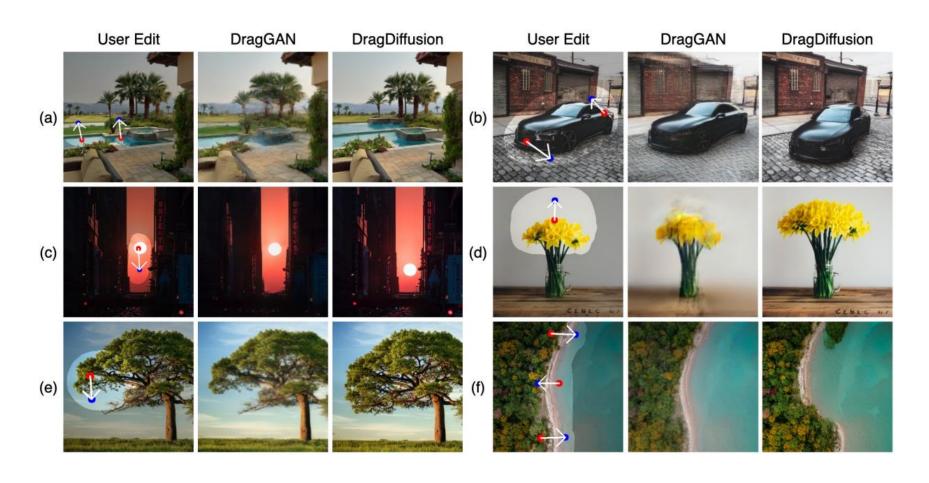
Improve consistency by referencing the correlated contents and texture of

original image

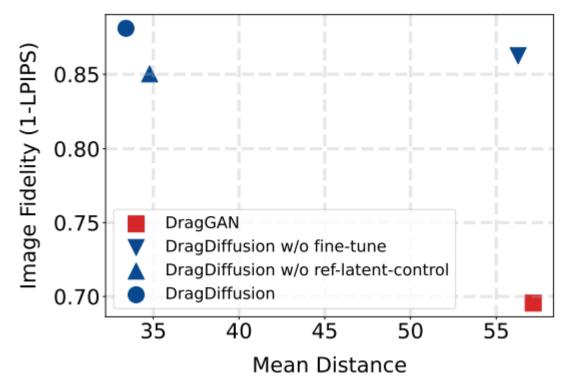


- Models: Stable Diffusion v1.5
- LoRA fine-tuning: 80 steps
- Evaluation
  - Image Fidelity ↑ (IF): quantifies the similarity between original and edited images
  - Mean Distance ↓ (MD): how well the approach moves the semantic contents to the target points

- All results are obtained under the same user edit
  - Measure the generality between GAN and Diffusion



- Comparison with DragGAN
  - Image Fidelity ↑ (IF): quantifies the similarity between original and edited images
  - Mean Distance ↓ (MD): how well the approach moves the semantic contents to the target points

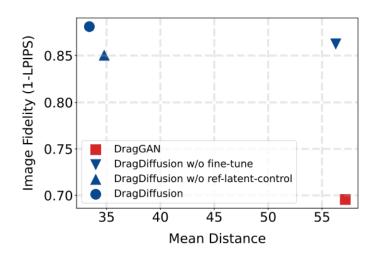


Ablation study

• Image Fidelity ↑ (IF): quantifies the similarity between original and edited images

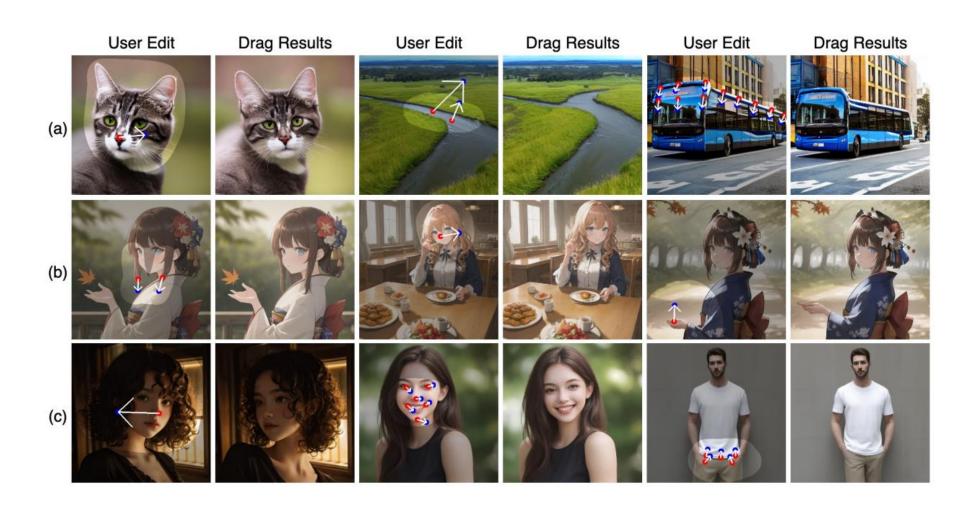
■ Mean Distance ↓ (MD): how well the approach moves the semantic contents

to the target points

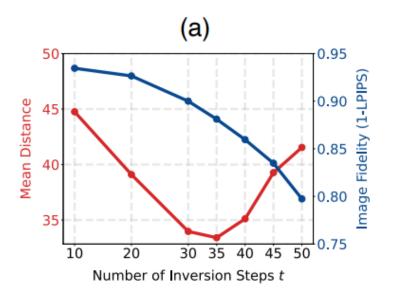


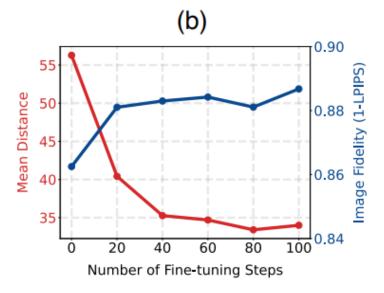


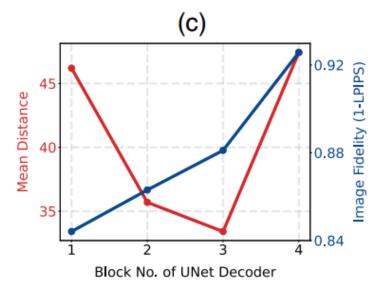
Show the generality of DragDiffusion



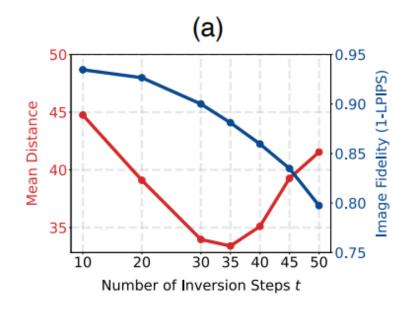
- Ablation study
  - (a): t = 35
  - (b): LoRA 80 steps
  - (c): which is better results to apply optimization loss

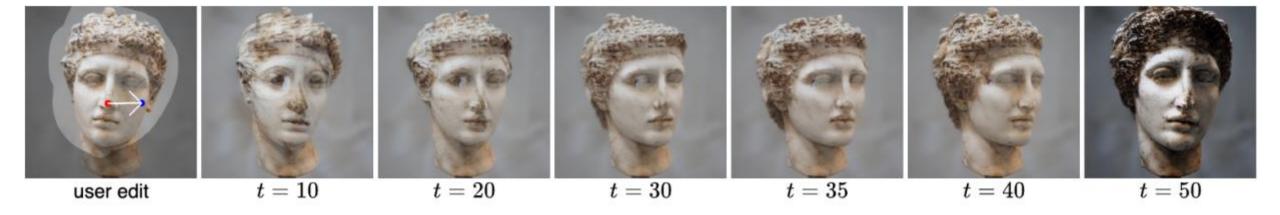




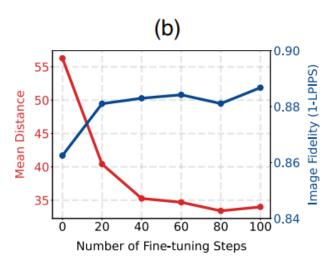


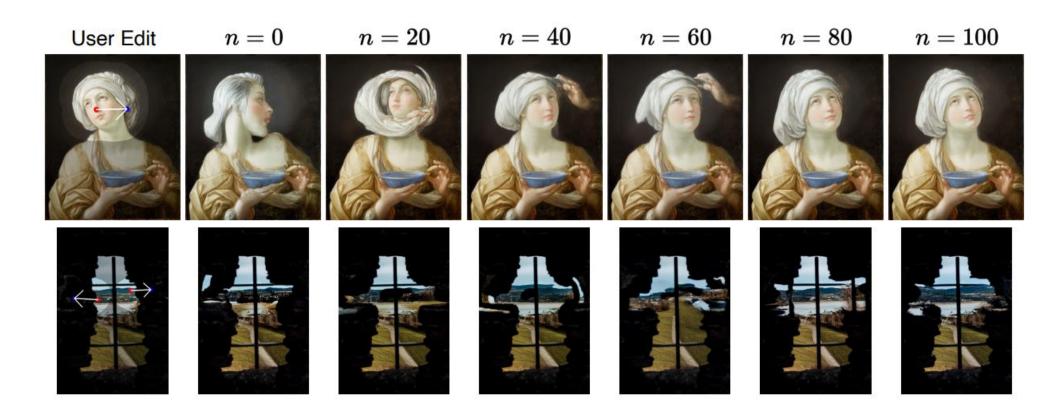
- Ablation study
  - Better results in  $t \in [30,40]$





- Ablation study
  - $n \ge 80$ : produce reasonable results without artifacts





- Ablation study
  - Conduct 3rd decoder block in UNet
  - 1st, 2nd: precise spatial control x
  - 4th: insufficient semantic and geometric information

