



北京大学
PEKING UNIVERSITY

Employing Diffusion Models for Low-Level Vision

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Visual Degradation

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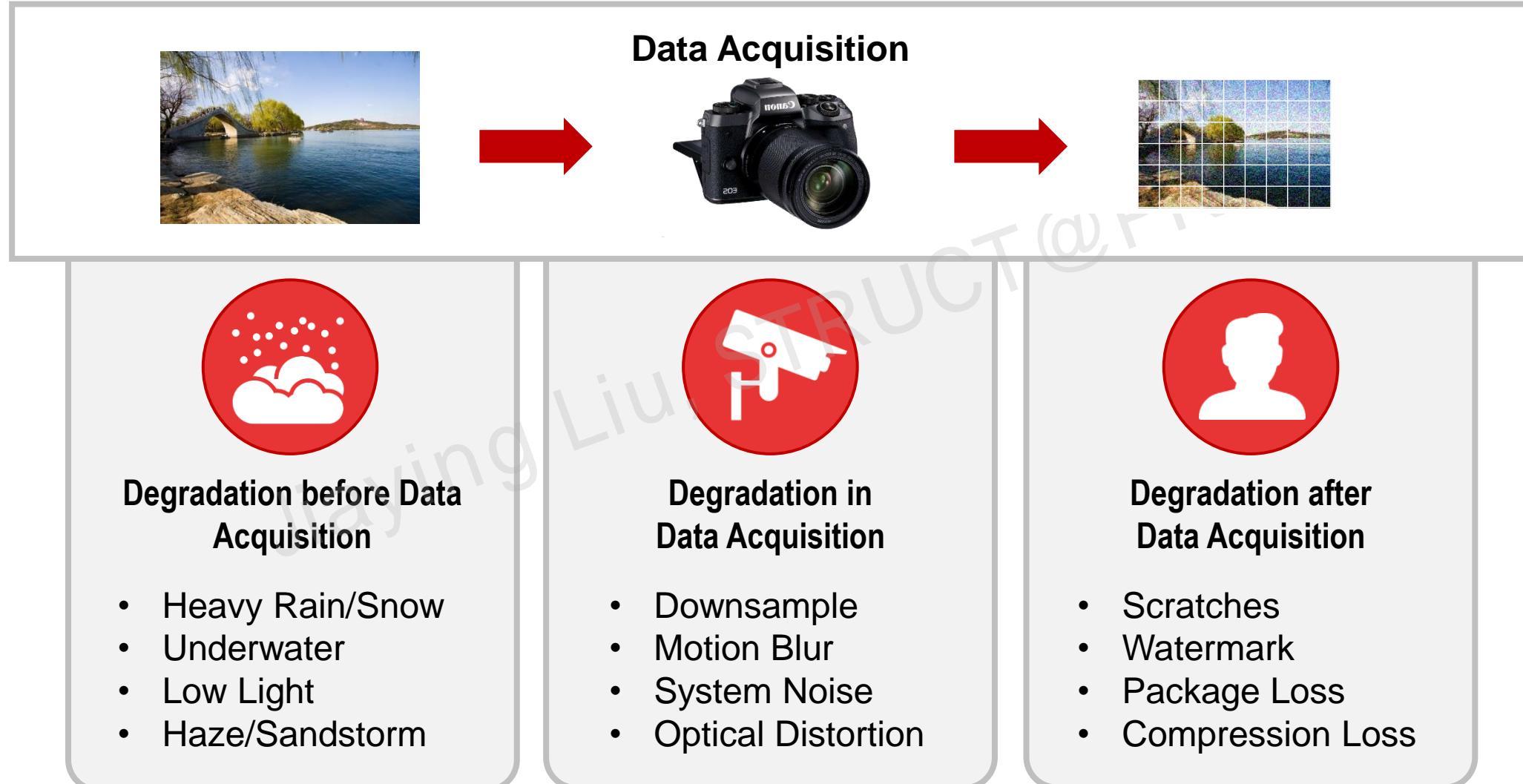


Image Restoration and Enhancement



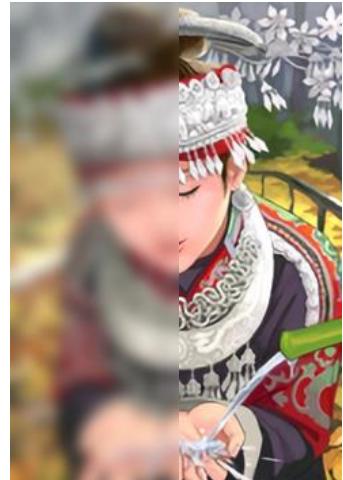
Underwater Enhancement



Dehazing



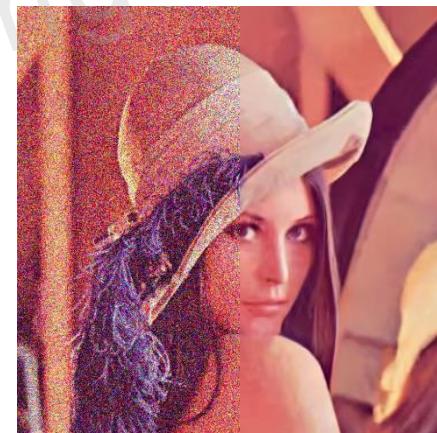
Text Removal



Super Resolution



Rain Streak Removal



Denoising



Rain Drop Removal

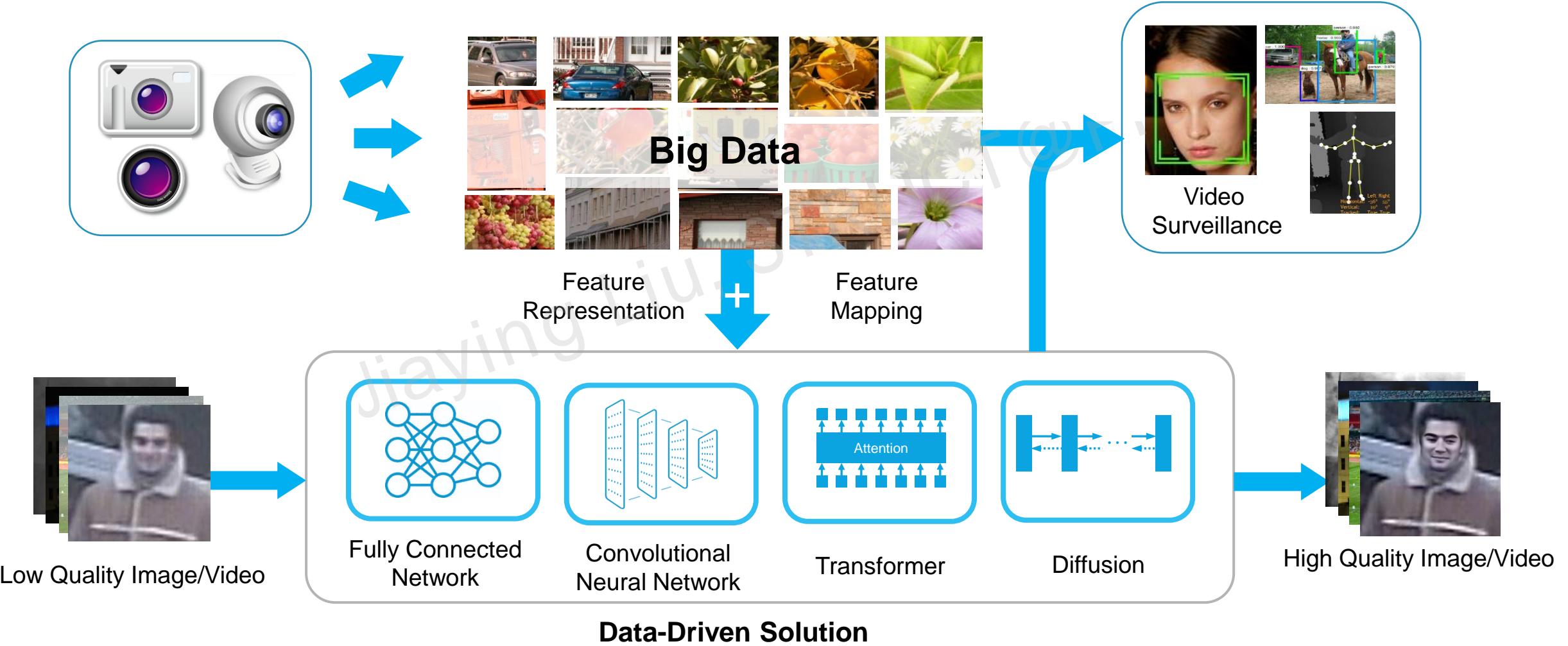


Low Light Enhancement

Image Restoration and Enhancement

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■ Data-Driven Solution



■ Diffusion for Image Restoration (IR)

Diffusion which creates visual details with high quality can play a role.

Text:

A gray sketch on paper of a Ferrari car, full car, pencil art

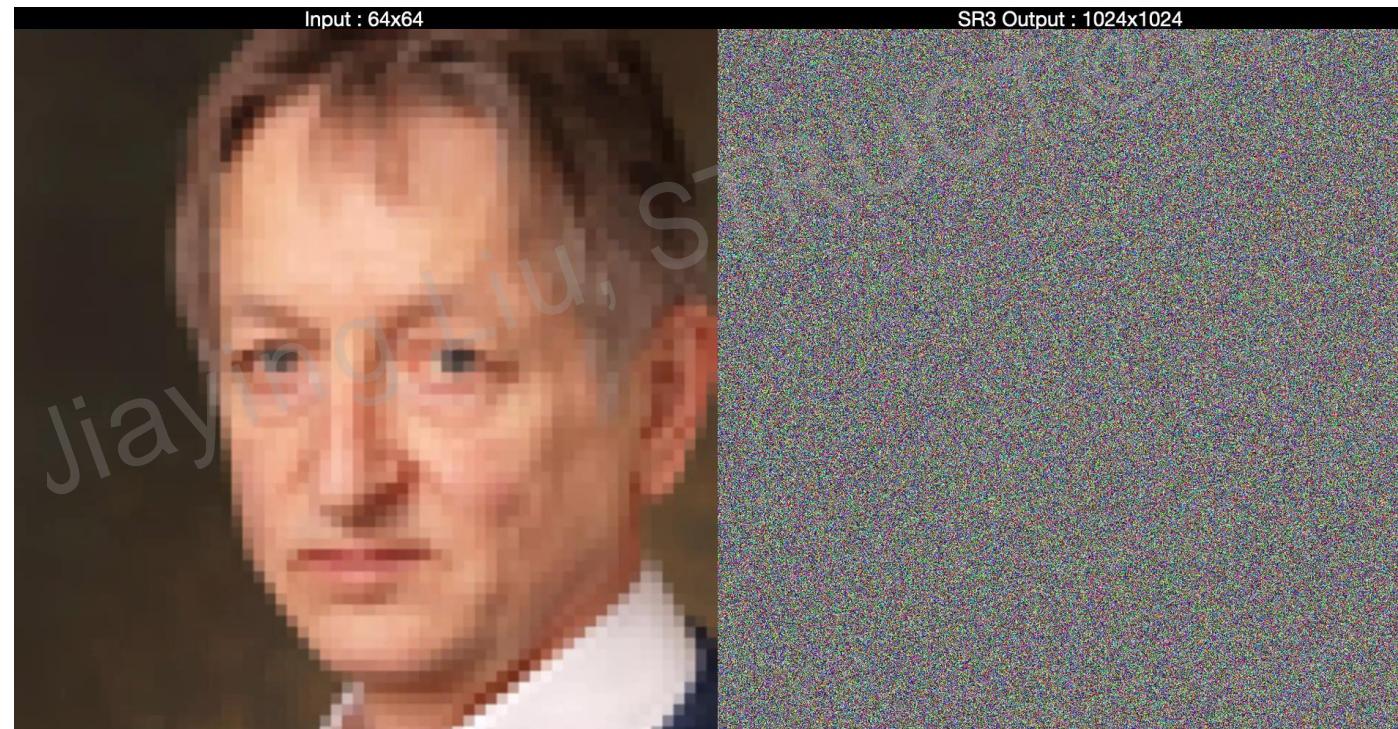


Image synthesis

[1] <https://prog.world/sherudim-under-the-hood-of-stable-diffusion/>

■ Diffusion for Image Restoration (IR)

Diffusion which creates visual details with high quality can play a role.



Restoration



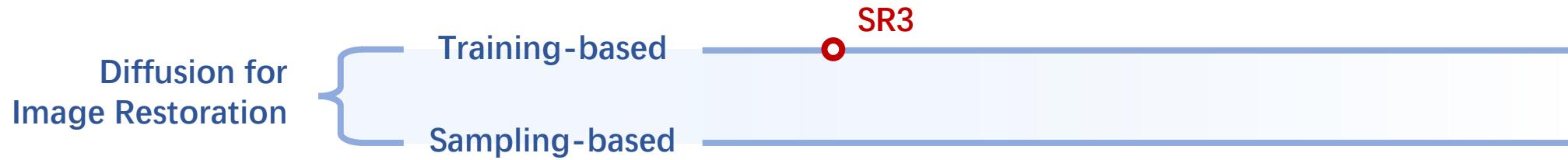
01

Diffusion for Image Restoration

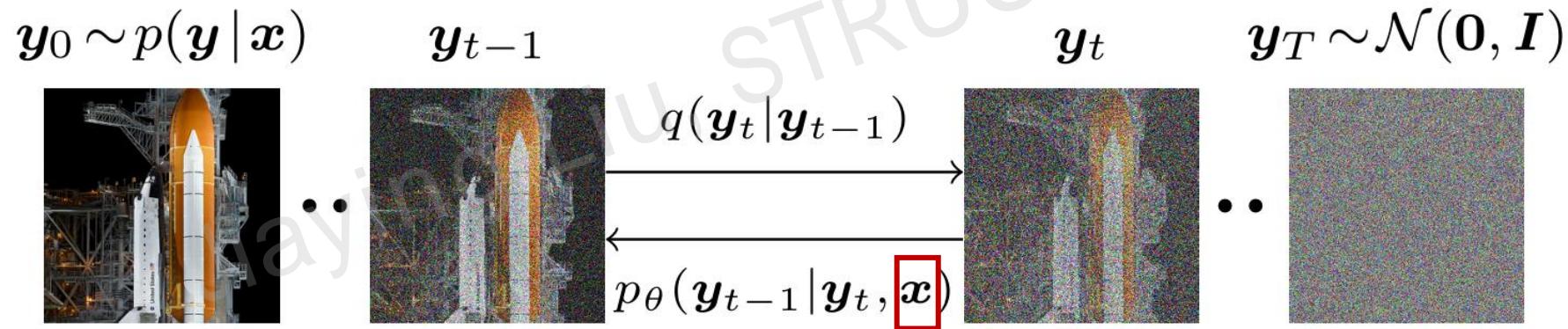
Diffusion for Image Restoration – SISR



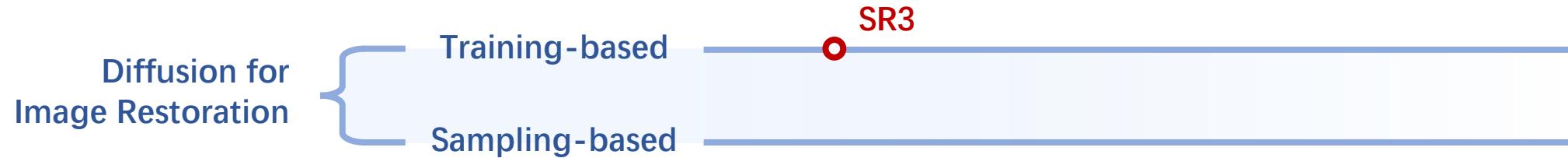
Diffusion for Image Restoration – SISR



■ Image Super-Resolution via Repeated Refinement (SR3)



- **Conditional** denoising diffusion model
- Forward diffusion process q , reverse inference process p
- Source image \mathbf{x} , target image \mathbf{y}_0



■ Image Super-Resolution via Repeated Refinement (SR3)

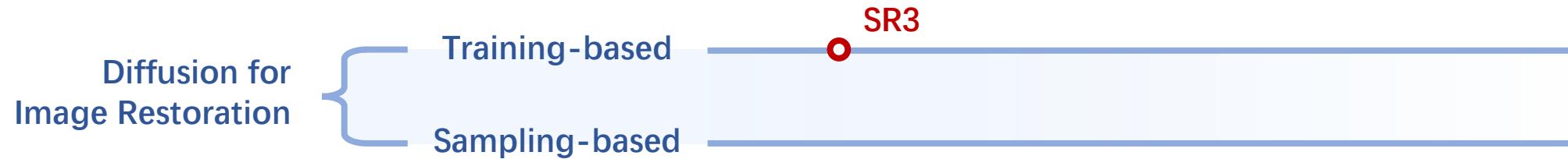
Conditioned reverse inference process

$$\begin{aligned} p_{\theta}(\mathbf{y}_{0:T} \mid \mathbf{x}) &= p(\mathbf{y}_T) \prod_{t=1}^T p_{\theta}(\mathbf{y}_{t-1} \mid \mathbf{y}_t, \mathbf{x}) \\ p(\mathbf{y}_T) &= \mathcal{N}(\mathbf{y}_T \mid \mathbf{0}, \mathbf{I}) \\ p_{\theta}(\mathbf{y}_{t-1} \mid \mathbf{y}_t, \mathbf{x}) &= \mathcal{N}(\mathbf{y}_{t-1} \mid \mu_{\theta}(\mathbf{x}, \mathbf{y}_t, \gamma_t), \sigma_t^2 \mathbf{I}) \end{aligned}$$

- The reverse inference process conditioned on source image \mathbf{x}
- $p_{\theta}(\mathbf{y}_{t-1} \mid \mathbf{y}_t, \mathbf{x})$ learned by the neural network

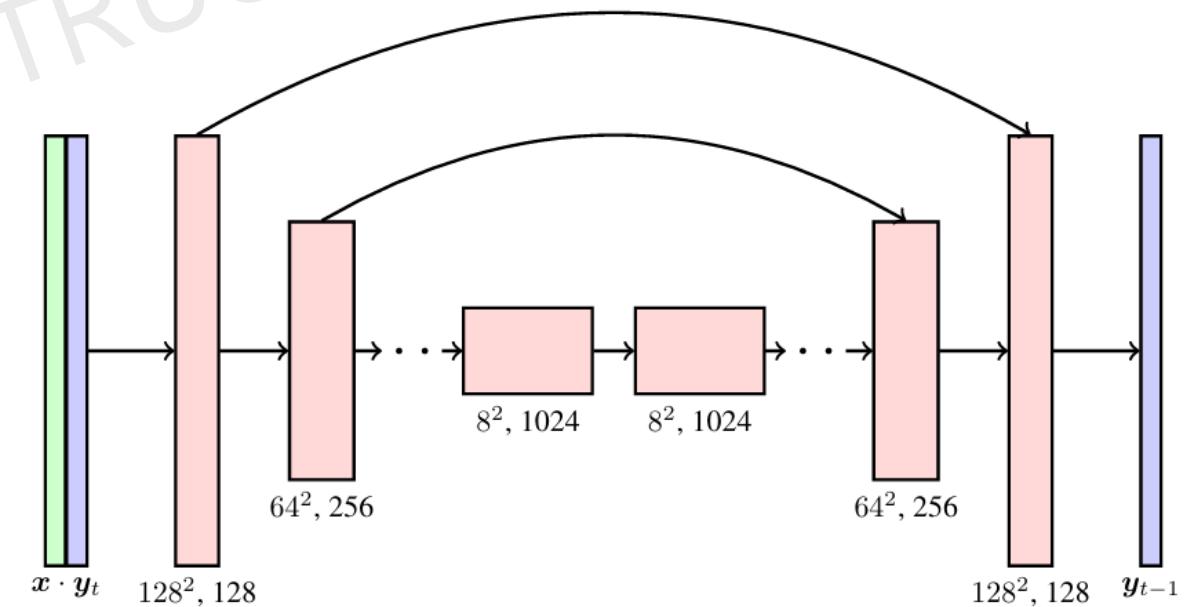
Diffusion for Image Restoration – SISR

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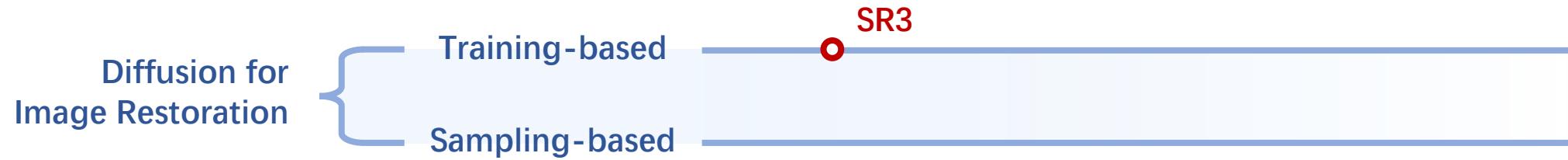
■ Image Super-Resolution via Repeated Refinement (SR3)

- U-Net architecture
- x interpolated to target high resolution
- x concatenate with noisy image y_t

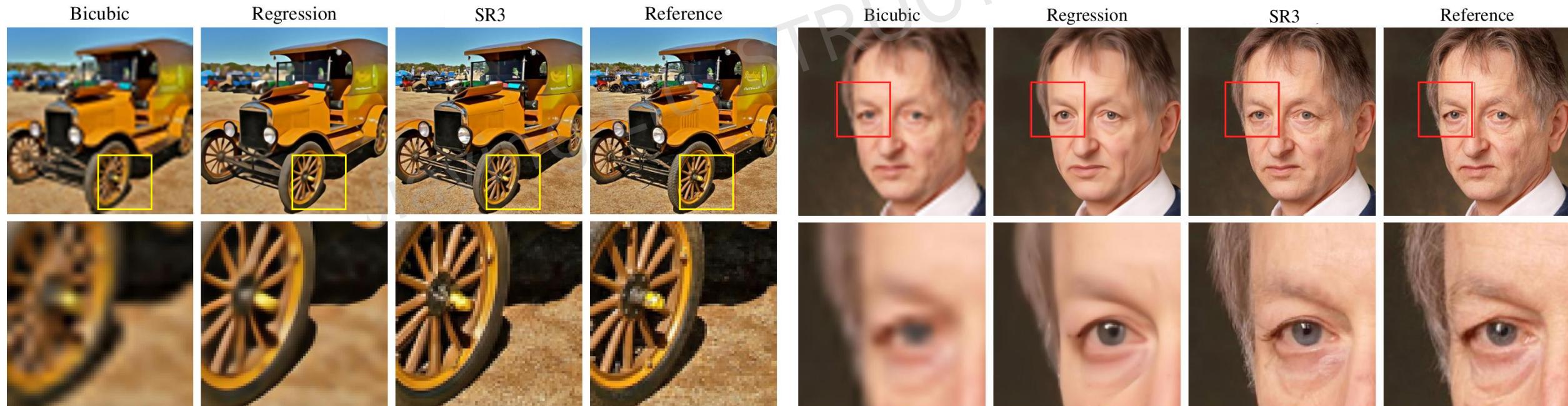


Diffusion for Image Restoration – SISR

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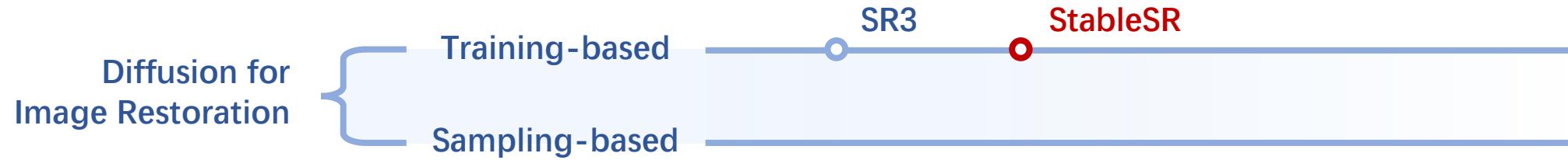


■ Image Super-Resolution via Repeated Refinement (SR3)



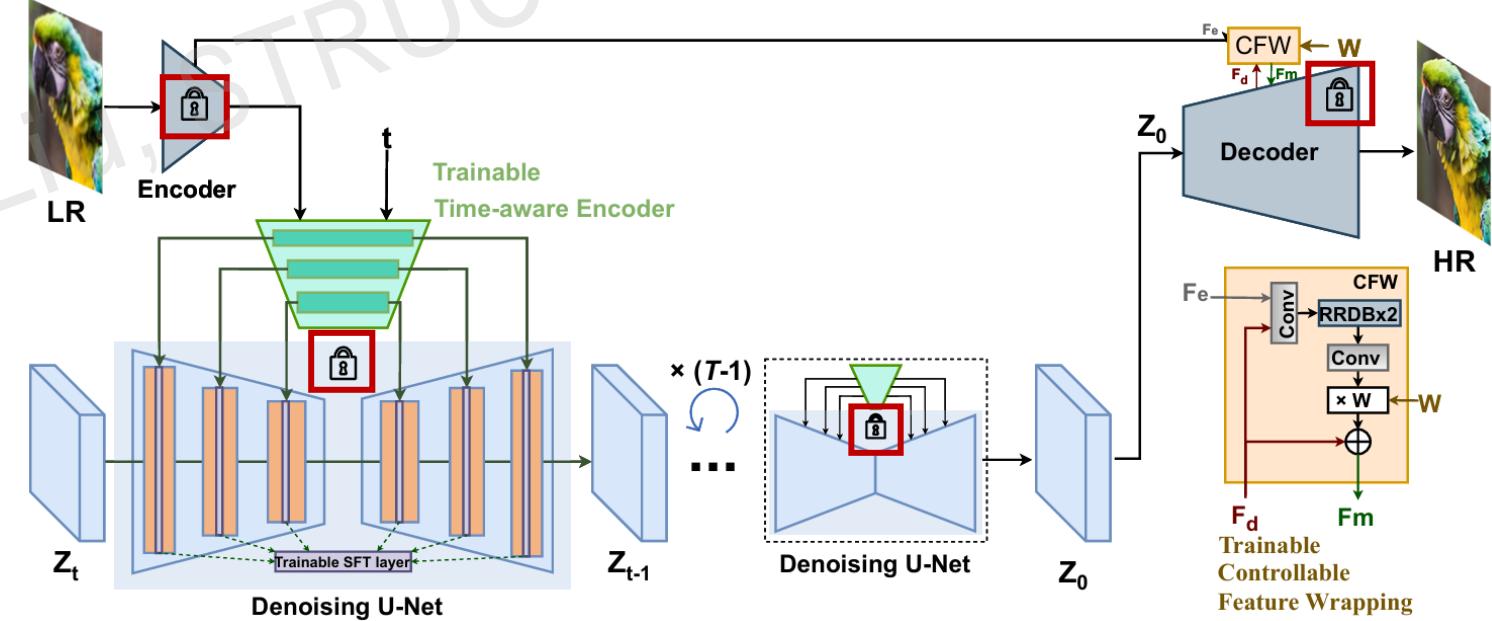
Diffusion for Image Restoration – SISR

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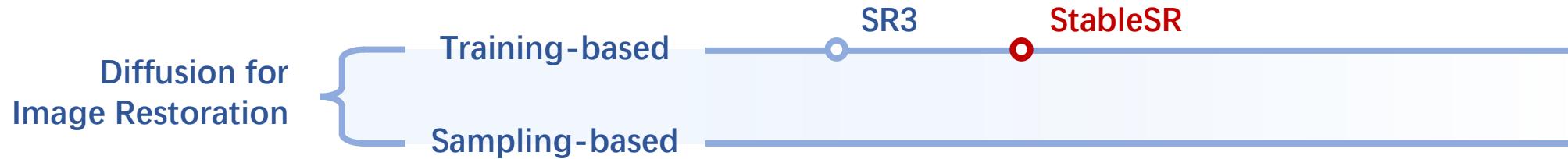
■ StableSR -- Image SR with diffusion prior

- Pretrained latent diffusion model for image synthesis



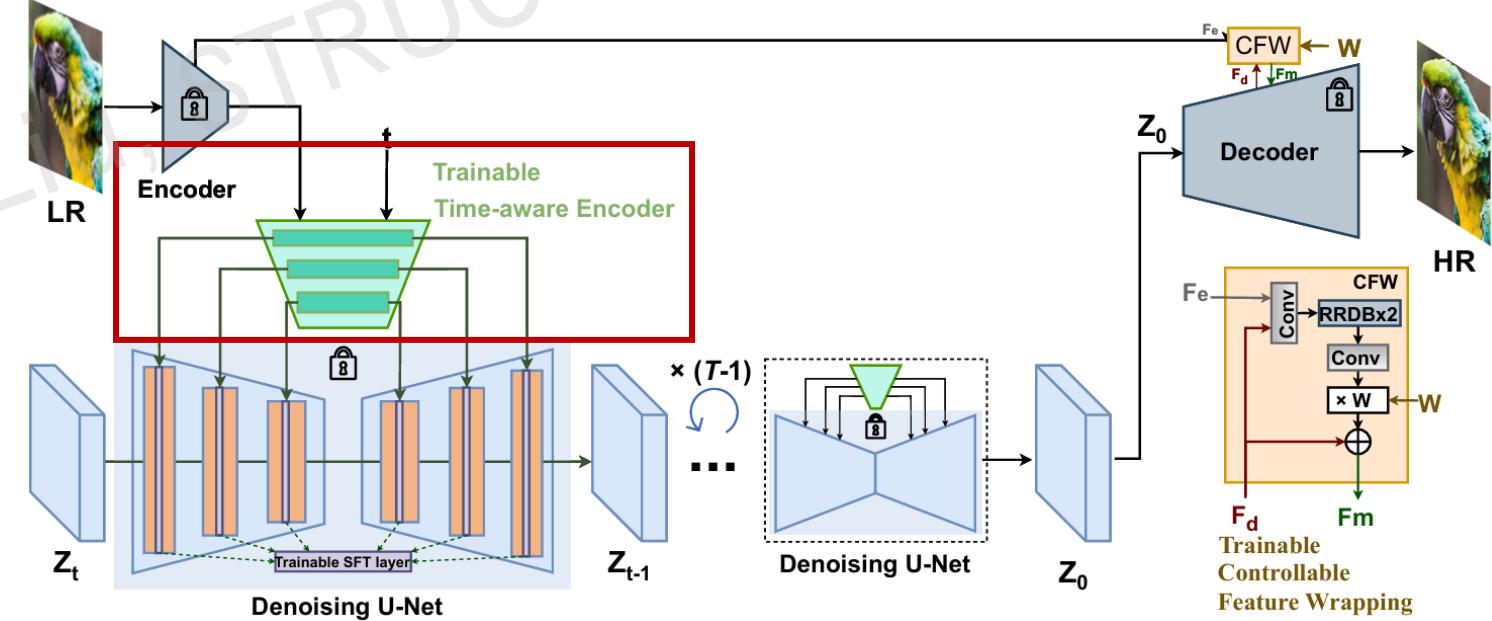
Diffusion for Image Restoration – SISR

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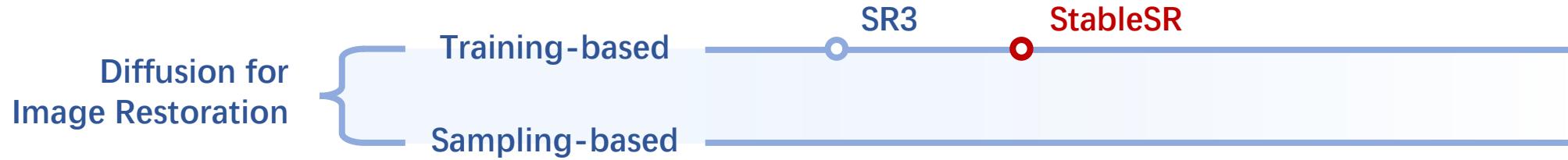
■ StableSR -- Image SR with diffusion prior

- Pretrained latent diffusion model for image synthesis
- Train time-aware encoder for SR



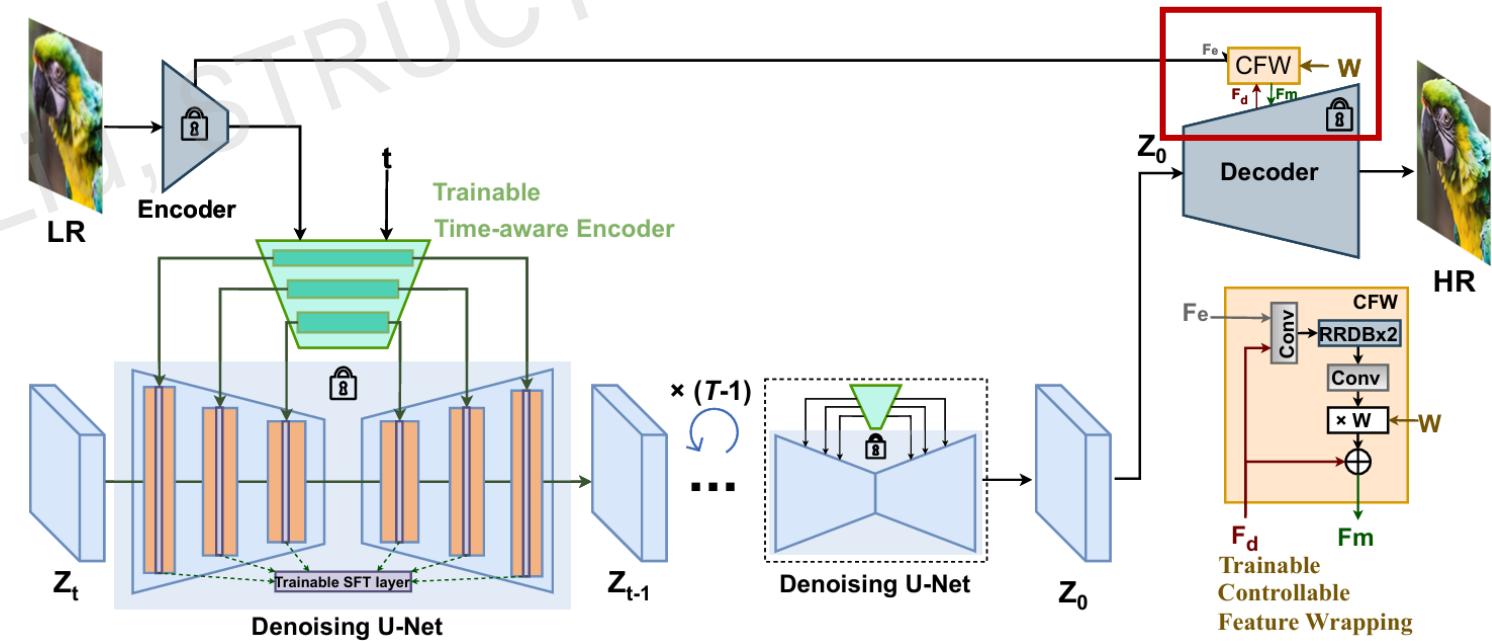
Diffusion for Image Restoration – SISR

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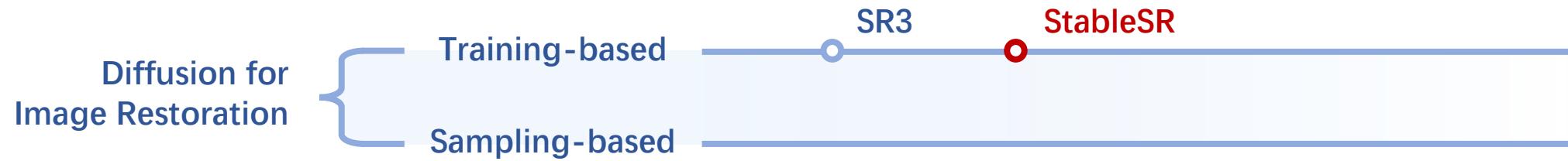
■ StableSR -- Image SR with diffusion prior

- Pretrained latent diffusion model for image synthesis
- Train time-aware encoder for SR
- Controllable feature wrapping to balance fidelity and perceptual quality

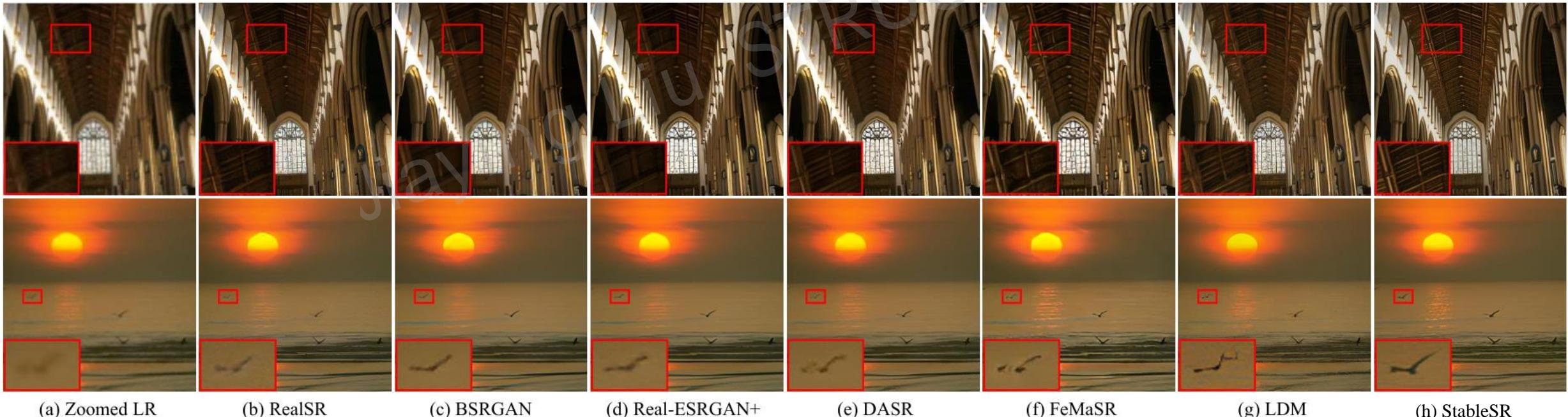


Diffusion for Image Restoration – SISR

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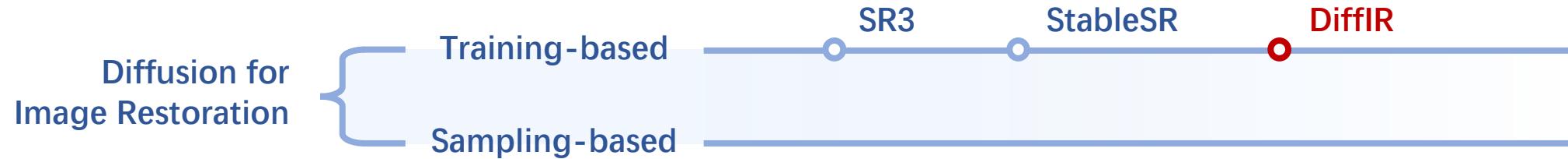


- StableSR -- Image SR with diffusion prior



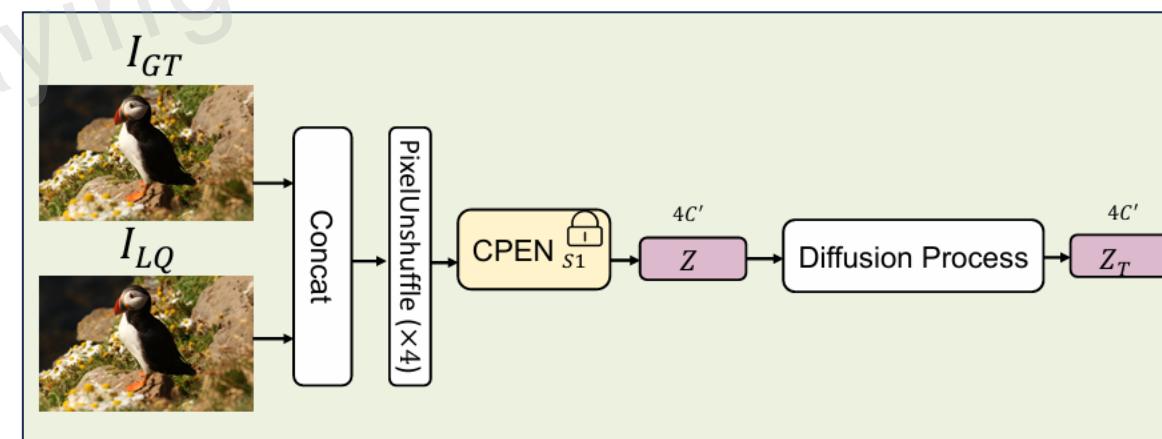
Diffusion for Image Restoration – SISR

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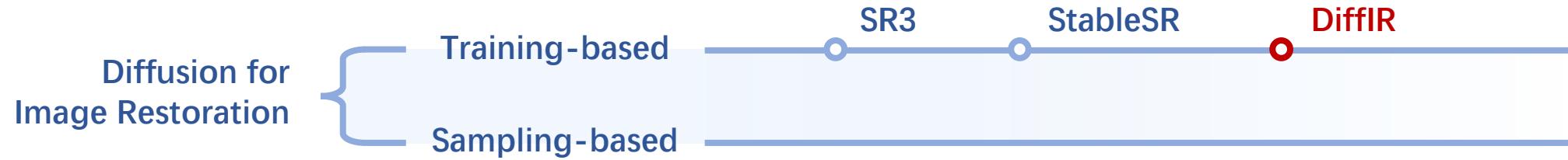
Diffusion Model for Image Restoration (DiffIR)

- Instead of generate target HR image with diffusion model, generate compact **IR prior**
- In **Stage 1**, IR prior Z is generated from ground truth and LR input, to guide restoration network



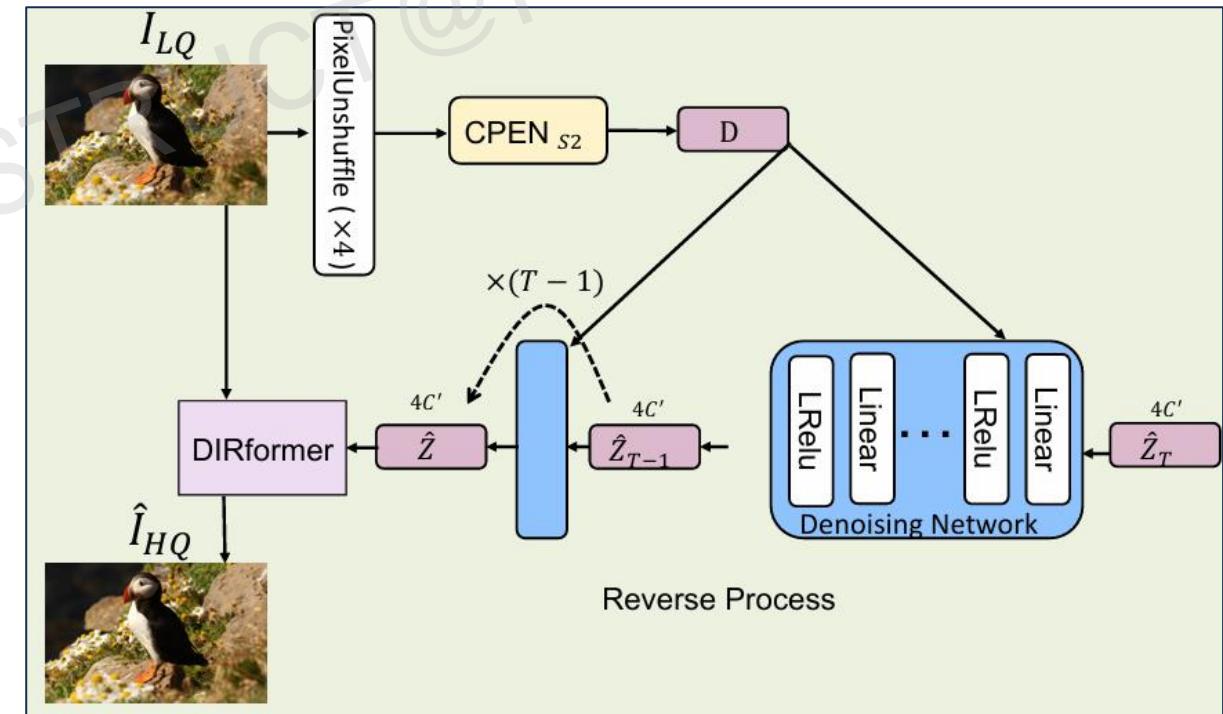
Diffusion for Image Restoration – SISR

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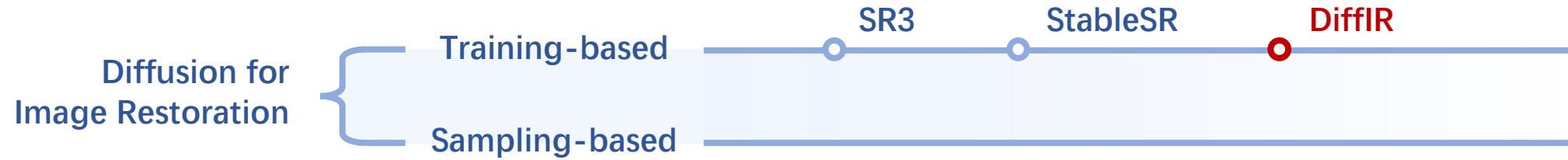
DiffIR

- In **Stage 2**, \hat{Z} is generated through conditional denoising diffusion model
- Guide restoration network DIRformer
- The generation of compact IR prior is faster and stabler



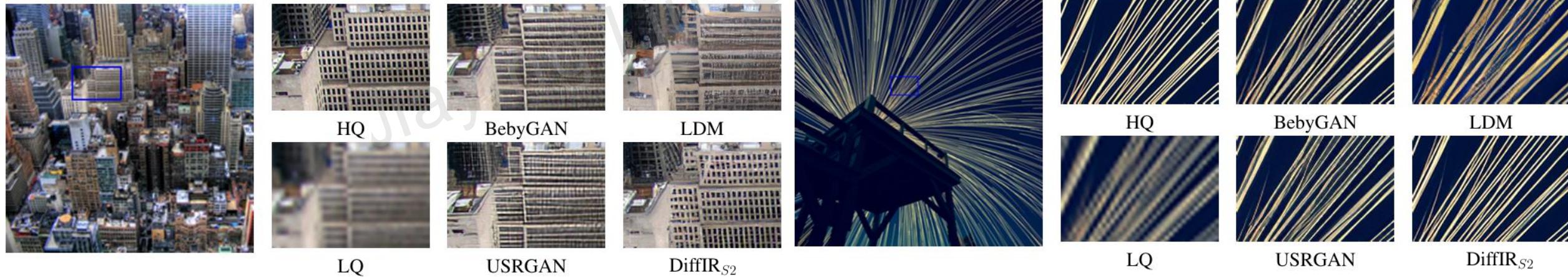
Diffusion for Image Restoration – SISR

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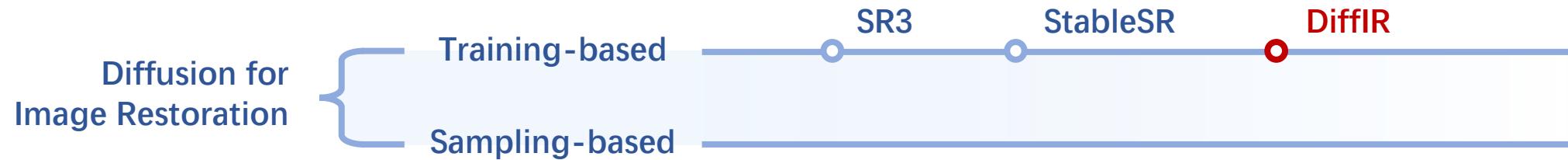
■ Diffusion model for image restoration (DiffIR)

■ Image super-resolution (4x)

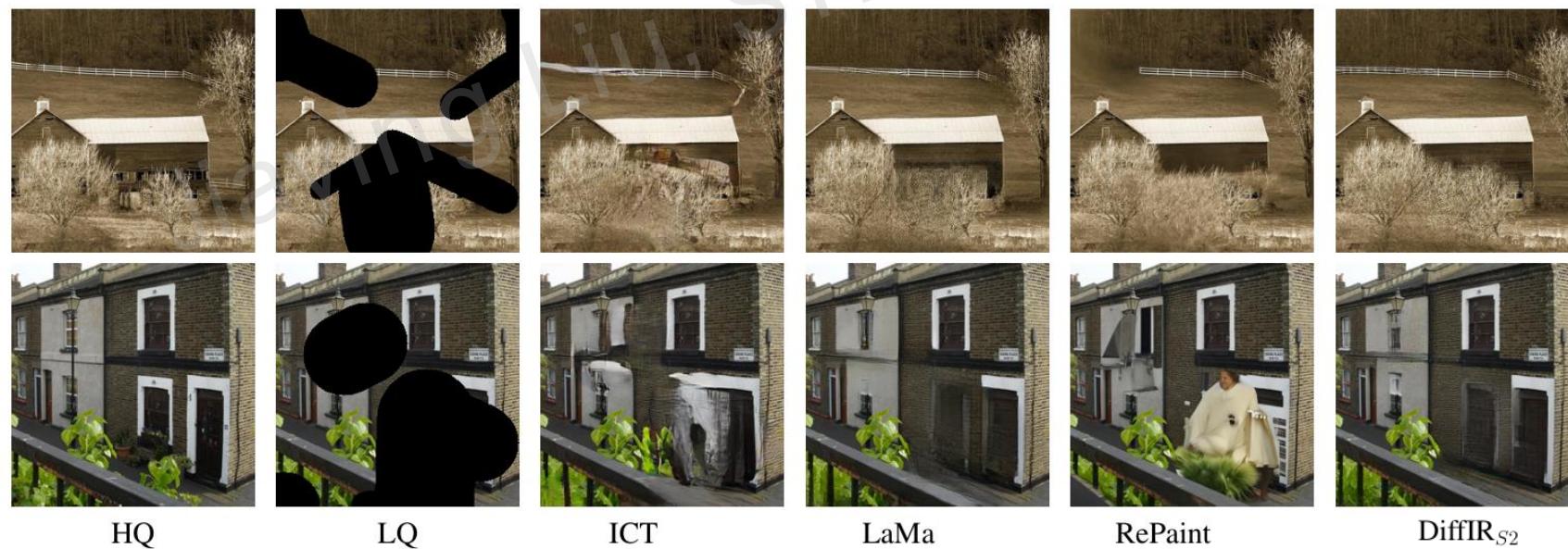


Diffusion for Image Restoration – SISR

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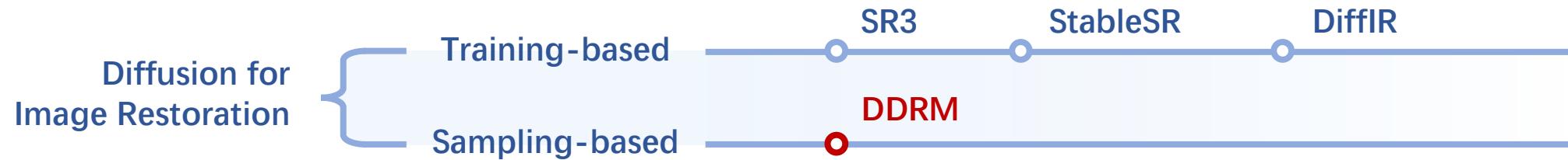


- Diffusion model for image restoration (DiffIR)
 - Inpainting task results



Diffusion for Image Restoration – SISR

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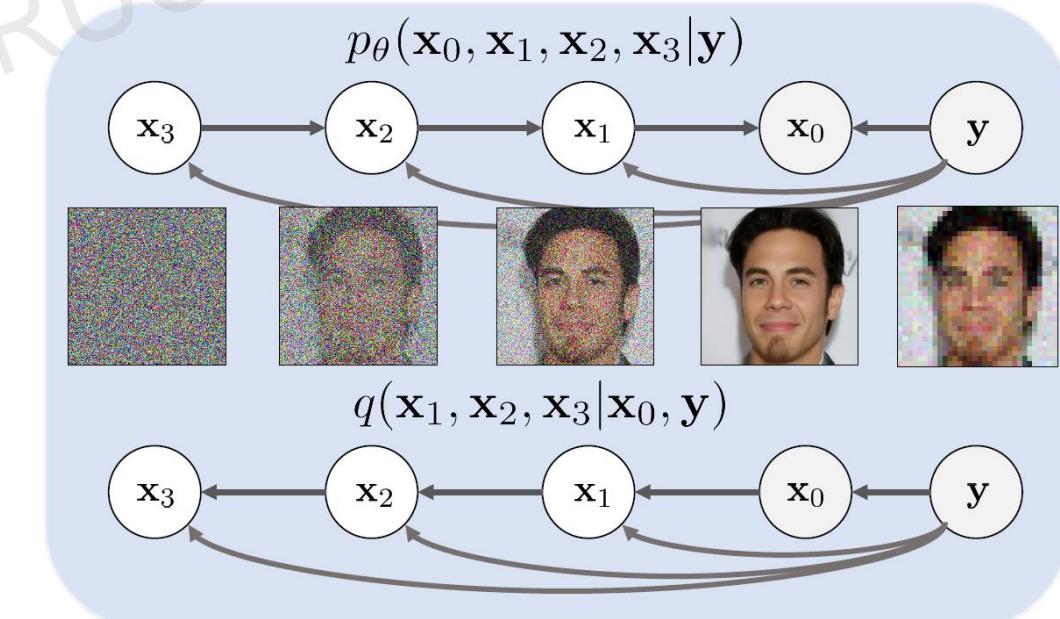


- Sampling restored images directly from pre-trained models

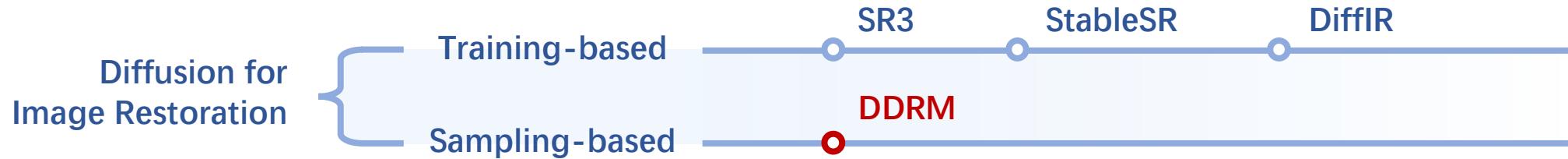
- Take the degraded image as the guidance

- Example:

Denoising Diffusion Restoration Models
(DDRM)



Denoising Diffusion Restoration Models



■ Denoising Diffusion Restoration Models

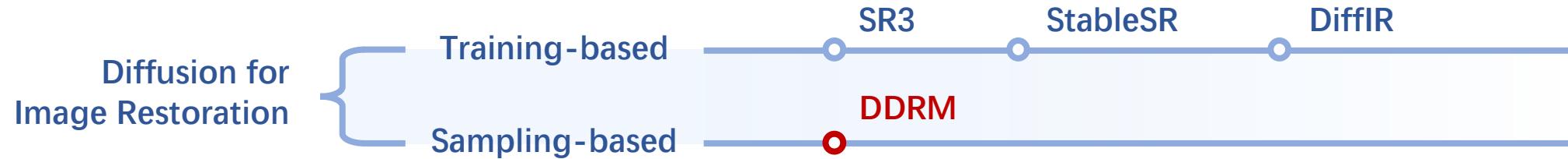
- A simplified degradation model in general:

$$\mathbf{y} = \mathbf{H}\mathbf{x} + \mathbf{z}, \quad \mathbf{z} \sim \mathcal{N}(0, \sigma_{\mathbf{y}}^2 \mathbf{I})$$

- The target is to find:

$$p_{\theta}(\mathbf{x}_{0:T} | \mathbf{y}) = p_{\theta}^{(T)}(\mathbf{x}_T | \mathbf{y}) \prod_{t=0}^{T-1} p_{\theta}^{(t)}(\mathbf{x}_t | \mathbf{x}_{t+1}, \mathbf{y})$$

where θ is the parameter set of a pre-trained diffusion model.



Denoising Diffusion Restoration Models

- The conditional distributions are demonstrated as:

Distribution of the initial noise

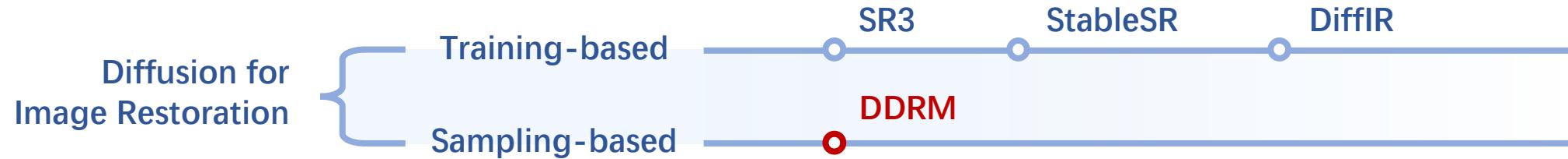
$$p_{\theta}^{(T)}(\bar{\mathbf{x}}_T^{(i)} | \mathbf{y}) = \begin{cases} \mathcal{N}(\bar{\mathbf{y}}^{(i)}, \sigma_T^2 - \frac{\sigma_y^2}{s_i^2}) & \text{if } s_i > 0 \\ \mathcal{N}(0, \sigma_T^2) & \text{if } s_i = 0 \end{cases}$$

$$p_{\theta}^{(t)}(\bar{\mathbf{x}}_t^{(i)} | \mathbf{x}_{t+1}, \mathbf{y}) = \begin{cases} \mathcal{N}(\bar{\mathbf{x}}_{\theta,t}^{(i)} + \sqrt{1 - \eta^2} \sigma_t \frac{\bar{\mathbf{x}}_{t+1}^{(i)} - \bar{\mathbf{x}}_{\theta,t}^{(i)}}{\sigma_{t+1}}, \eta^2 \sigma_t^2) & \text{if } s_i = 0 \\ \mathcal{N}(\bar{\mathbf{x}}_{\theta,t}^{(i)} + \sqrt{1 - \eta^2} \sigma_t \frac{\bar{\mathbf{y}}^{(i)} - \bar{\mathbf{x}}_{\theta,t}^{(i)}}{\sigma_y / s_i}, \eta^2 \sigma_t^2) & \text{if } \sigma_t < \frac{\sigma_y}{s_i} \\ \mathcal{N}((1 - \eta_b) \bar{\mathbf{x}}_{\theta,t}^{(i)} + \eta_b \bar{\mathbf{y}}^{(i)}, \sigma_t^2 - \frac{\sigma_y^2}{s_i^2} \eta_b^2) & \text{if } \sigma_t \geq \frac{\sigma_y}{s_i} \end{cases}$$

Perform SVD $H = U\Sigma V^\top$

Correct Sampling Distribution

where $\bar{\mathbf{x}}_t$ is the set of vectors in spectral space: $\bar{\mathbf{x}}_t = V^\top \mathbf{x}_t$. Requires known degradation.



Denoising Diffusion Restoration Models

- The conditional distributions are demonstrated as:

$$p_{\theta}^{(T)}(\bar{\mathbf{x}}_T^{(i)} | \mathbf{y}) = \begin{cases} \mathcal{N}(\bar{\mathbf{y}}^{(i)}, \sigma_T^2 - \frac{\sigma_y^2}{s_i^2}) & \text{if } s_i > 0 \\ \mathcal{N}(0, \sigma_T^2) & \text{if } s_i = 0 \end{cases}$$

Distribution of one single denoise step

$$p_{\theta}^{(t)}(\bar{\mathbf{x}}_t^{(i)} | \mathbf{x}_{t+1}, \mathbf{y}) = \begin{cases} \mathcal{N}(\bar{\mathbf{x}}_{\theta,t}^{(i)} + \sqrt{1 - \eta^2} \sigma_t \frac{\bar{\mathbf{x}}_{t+1}^{(i)} - \bar{\mathbf{x}}_{\theta,t}^{(i)}}{\sigma_{t+1}}, \eta^2 \sigma_t^2) & \text{if } s_i = 0 \\ \mathcal{N}(\bar{\mathbf{x}}_{\theta,t}^{(i)} + \sqrt{1 - \eta^2} \sigma_t \frac{\bar{\mathbf{y}}^{(i)} - \bar{\mathbf{x}}_{\theta,t}^{(i)}}{\sigma_y / s_i}, \eta^2 \sigma_t^2) & \text{if } \sigma_t < \frac{\sigma_y}{s_i} \\ \mathcal{N}((1 - \eta_b) \bar{\mathbf{x}}_{\theta,t}^{(i)} + \eta_b \bar{\mathbf{y}}^{(i)}, \sigma_t^2 - \frac{\sigma_y^2}{s_i^2} \eta_b^2) & \text{if } \sigma_t \geq \frac{\sigma_y}{s_i} \end{cases}$$

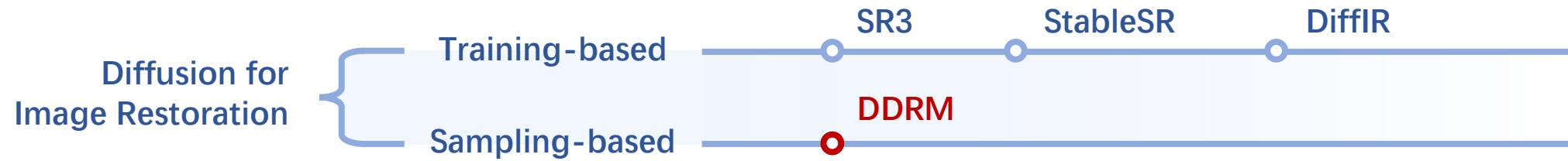
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Diffusion for Image Restoration – SISR

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Denoising Diffusion Restoration Models

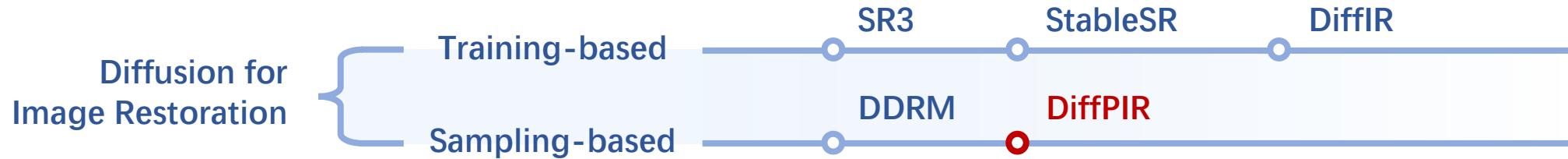


(a) Inpainting results on cat images.

(b) Deblurring results ($\sigma_y = 0.05$) on bedroom images.

Diffusion for Image Restoration – SISR

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■ Denoising Diffusion Models for Plug-and-Play Image Restoration (DiffPIR)



Employ a refinement during sampling process:

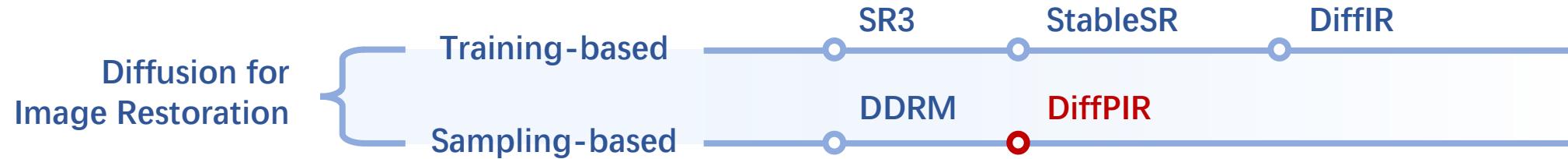
$$\hat{\mathbf{x}}_0^{(t)} = \arg \min_{\mathbf{x}} \|\mathbf{y} - \mathcal{H}(\mathbf{x})\|^2 + \rho_t \|\mathbf{x} - \mathbf{x}_0^{(t)}\|^2$$

Fidelity to degraded image

where $\mathbf{x}_0^{(t)}$ is the predicted noise-free image at timestep t . Requires known degradation.

Diffusion for Image Restoration – SISR

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■ Denoising Diffusion Models for Plug-and-Play Image Restoration



Employ a refinement during sampling process:

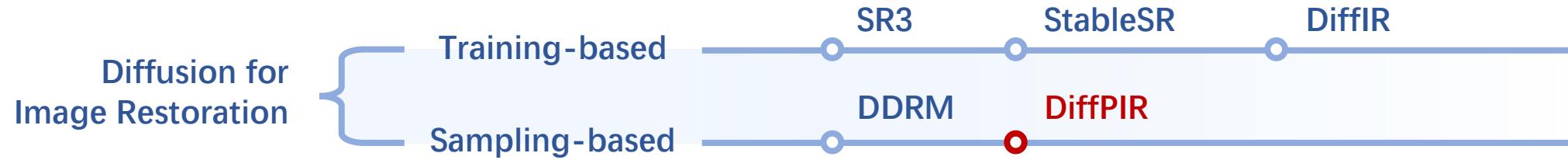
$$\hat{\mathbf{x}}_0^{(t)} = \arg \min_{\mathbf{x}} \|\mathbf{y} - \mathcal{H}(\mathbf{x})\|^2 + \rho_t \|\mathbf{x} - \mathbf{x}_0^{(t)}\|^2$$

Regularization of generated image

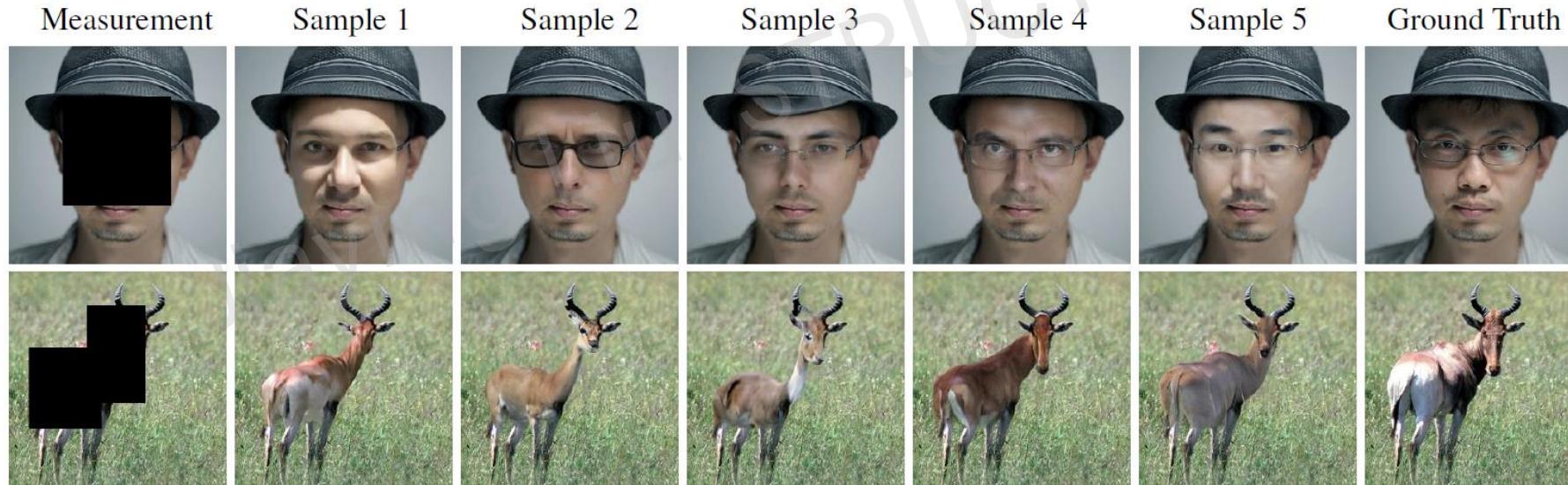
where $\mathbf{x}_0^{(t)}$ is the predicted noise-free image at timestep t . Requires known degradation.

Diffusion for Image Restoration – SISR

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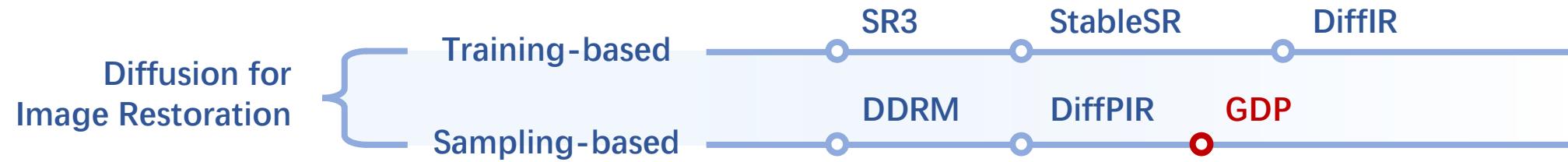
■ Denoising Diffusion Models for Plug-and-Play Image Restoration



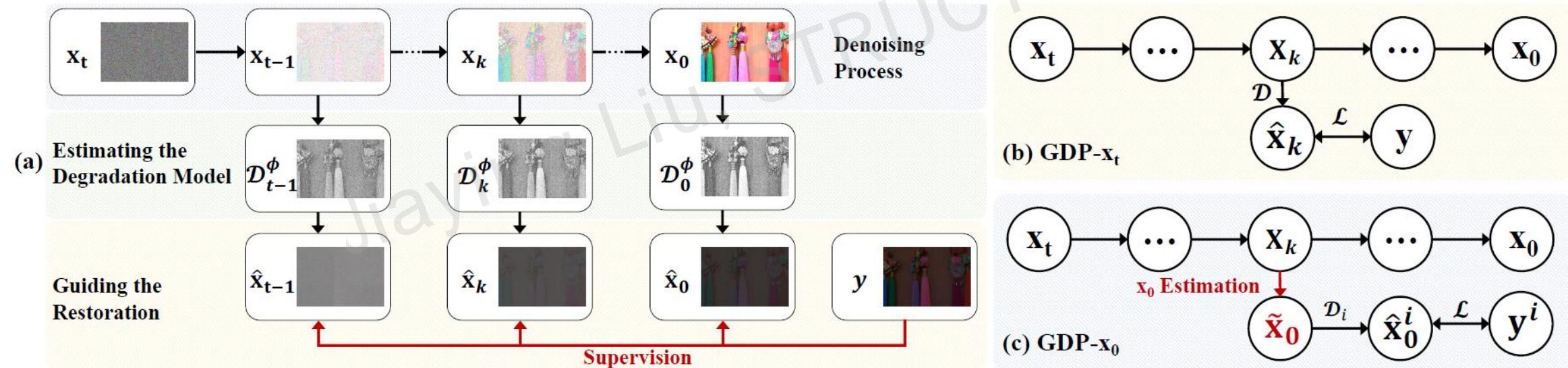
Diverse and reasonable results.

Diffusion for Image Restoration – SISR

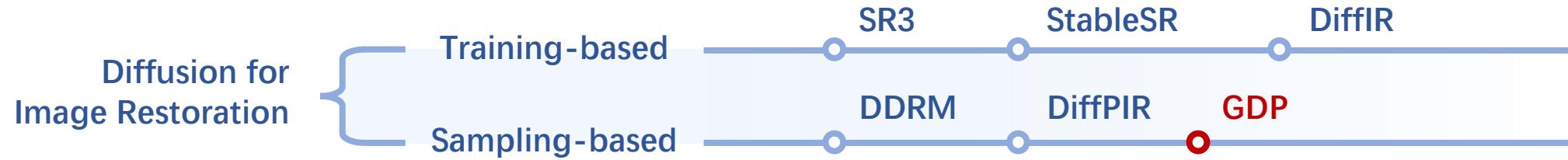
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■ Generative Diffusion Prior (GDP)



Learn the degradation model during the sampling process.



■ Generative Diffusion Prior (GDP)

At the t -th step, the degradation model will be updated via:

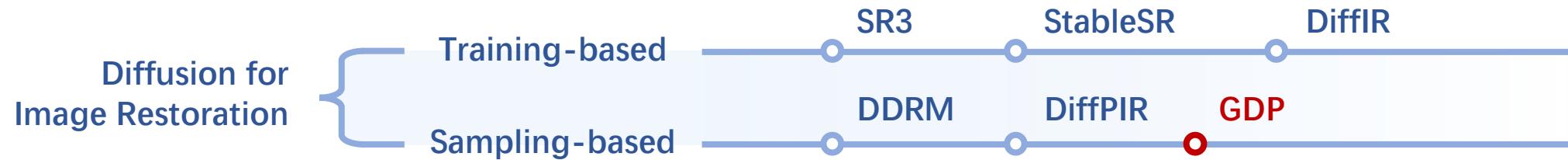
$$\begin{aligned}\mathcal{L}_{\phi, \tilde{x}_0}^{total} &= \mathcal{L}(y, \boxed{\mathcal{D}_\phi}(\tilde{x}_0)) + \mathcal{Q}(\tilde{x}_0) \\ \phi &\leftarrow \phi - l \nabla_\phi \mathcal{L}_{\phi, \tilde{x}_0}^{total}\end{aligned}$$

The noisy image at $(t-1)$ -th step will be sampled with guidance provided by the loss:

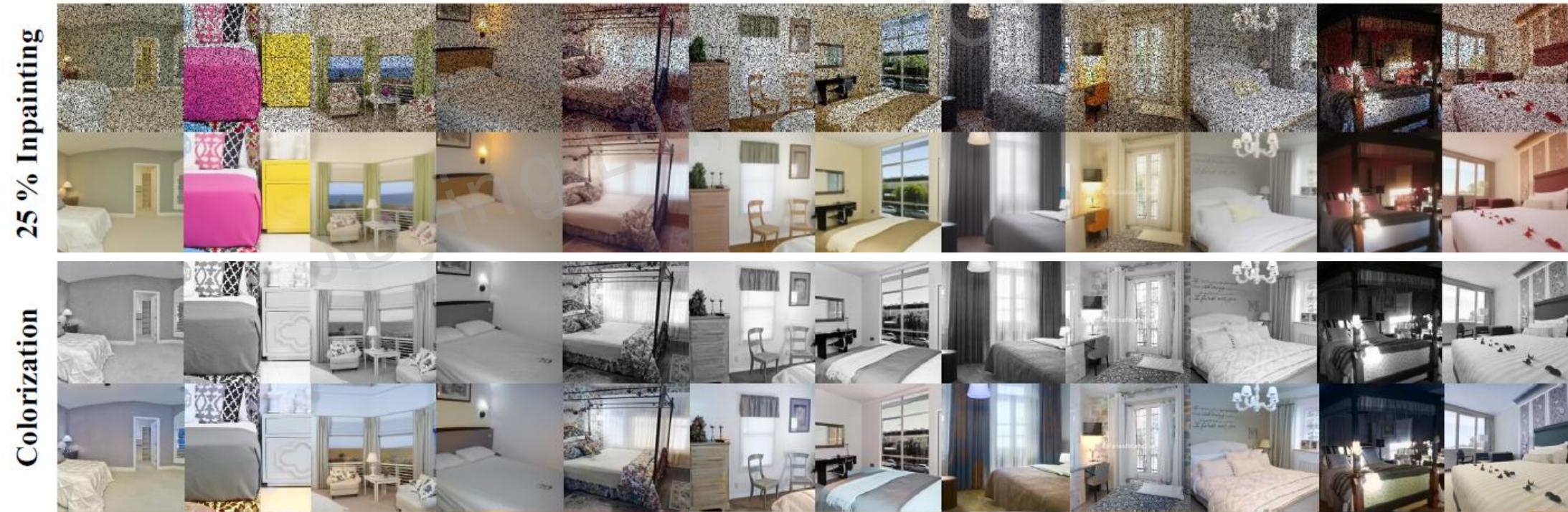
$$\text{Sample } \mathbf{x}_{t-1} \text{ by } \mathcal{N} \left(\mu + s \nabla_{\tilde{x}_0} \mathcal{L}_{\phi, \tilde{x}_0}^{total}, \Sigma \right)$$

Diffusion for Image Restoration – SISR

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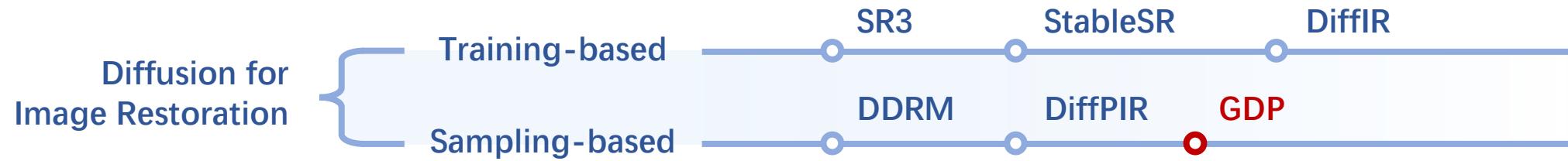


■ Generative Diffusion Prior (GDP)

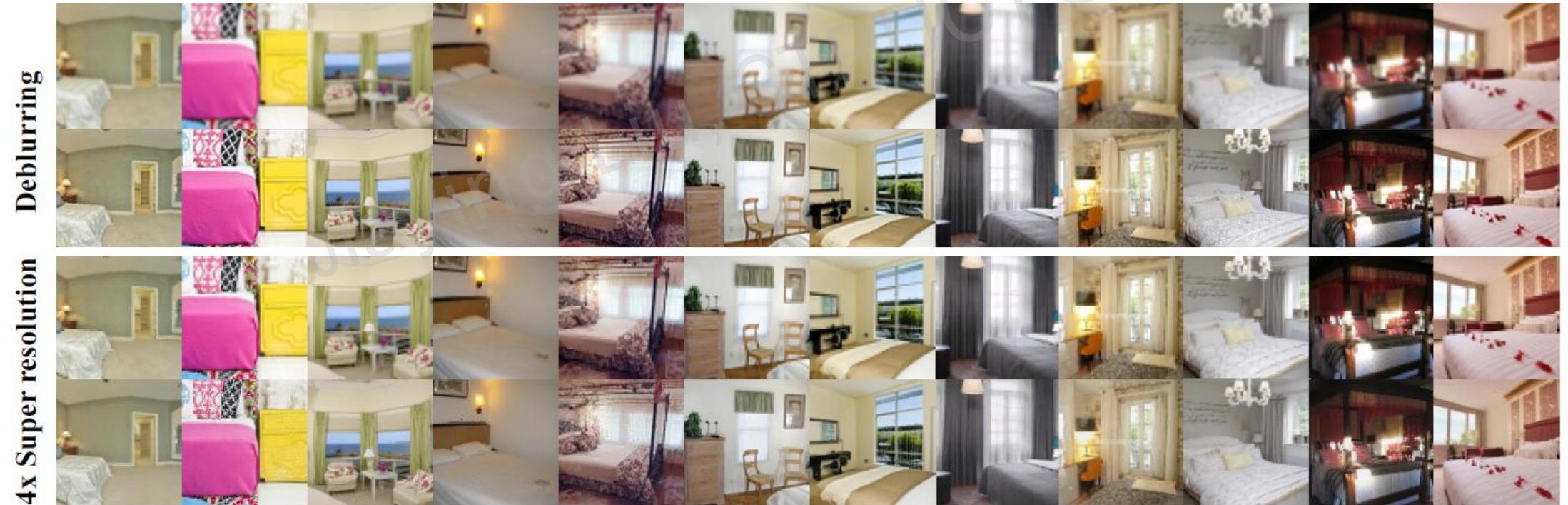


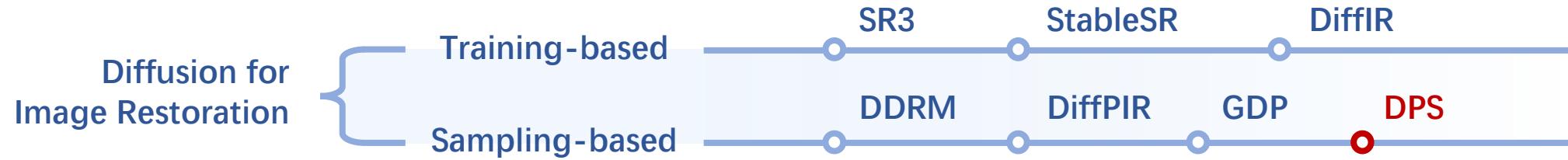
Diffusion for Image Restoration – SISR

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■ Generative Diffusion Prior (GDP)

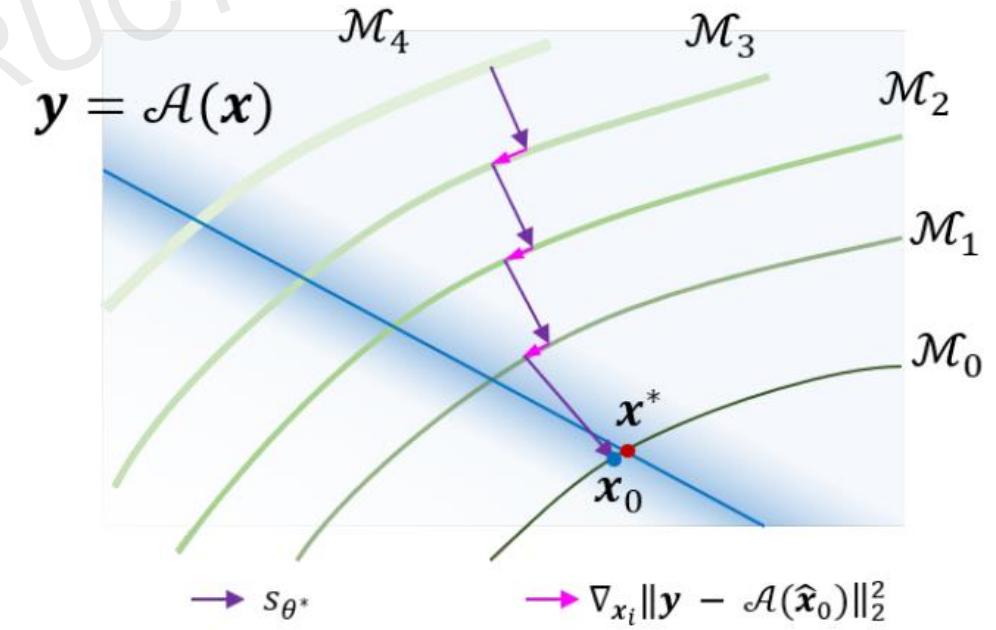


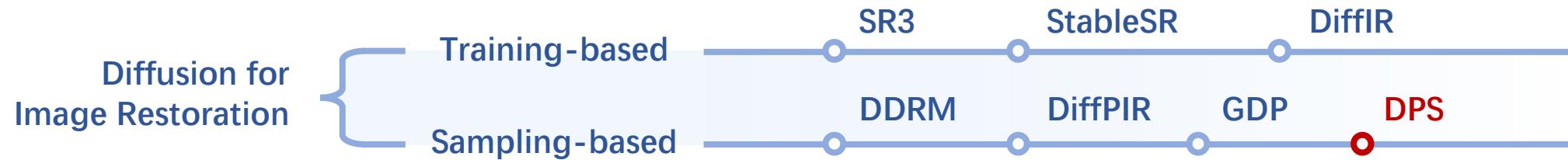


■ Diffusion Posterior Sampling (DPS)

refine the noisy state at sampling process via:

$$\mathbf{x}_{i-1} \leftarrow \mathbf{x}'_{i-1} - \zeta_i \nabla_{\mathbf{x}_i} \|\mathbf{y} - \mathcal{A}(\hat{\mathbf{x}}_0)\|_2^2$$



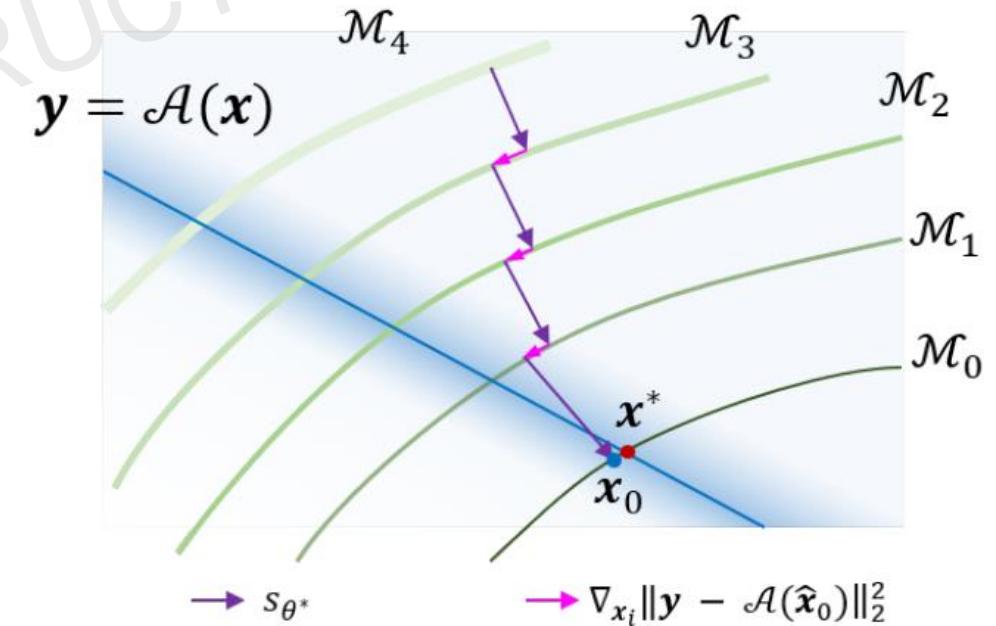


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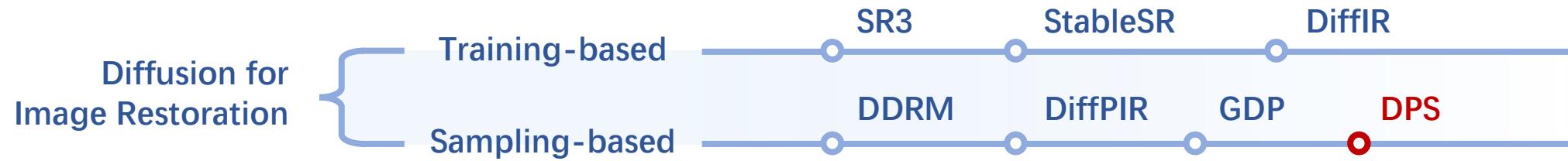
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Correction Item

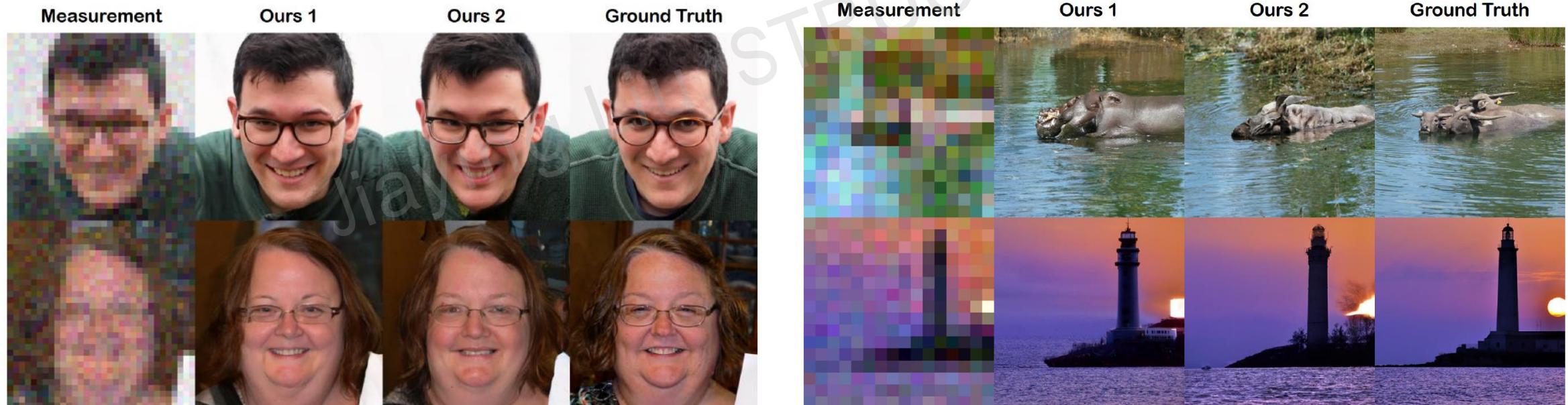


Diffusion for Image Restoration – SISR

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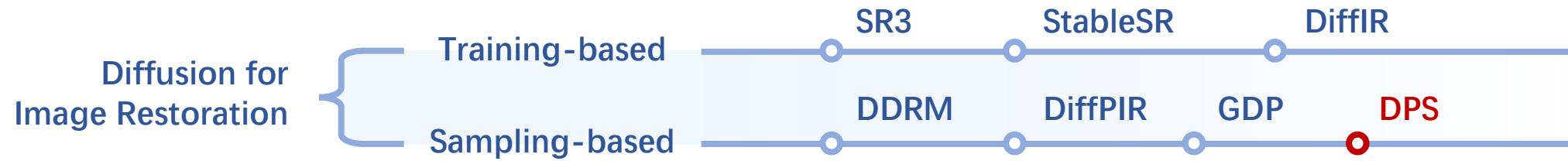


Diffusion Posterior Sampling (DPS)

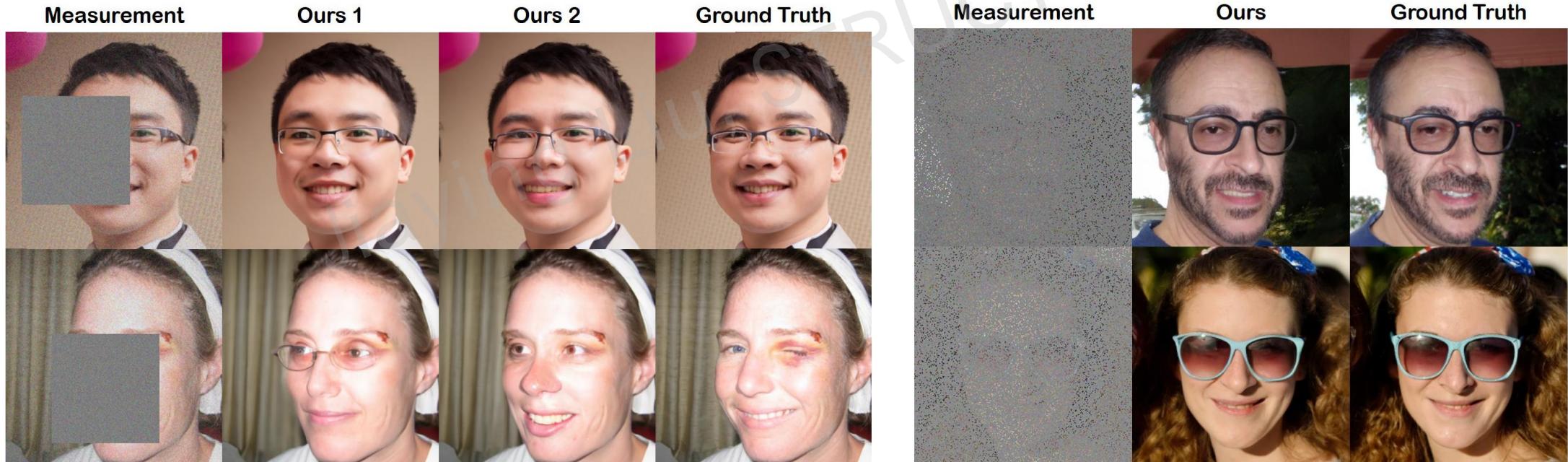


Diffusion for Image Restoration – SISR

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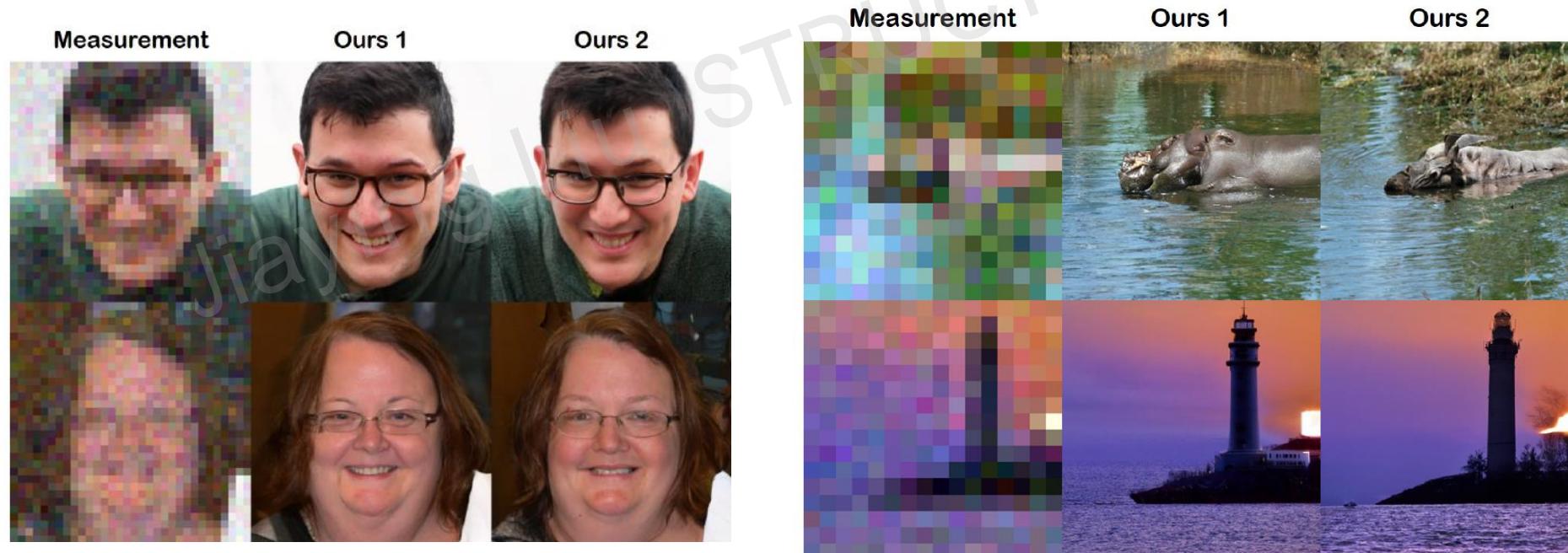
Diffusion Posterior Sampling (DPS)



■ Problem of existing works:

There are two characteristics of diffusion-based IR models:

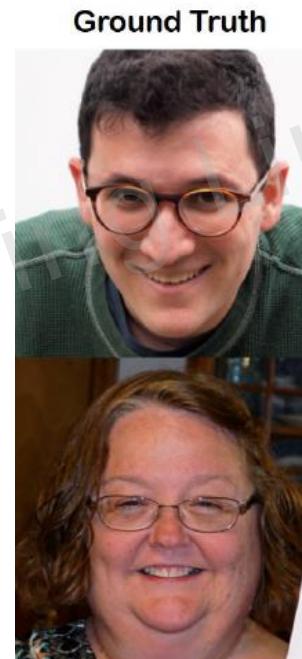
- Sampling **instability** due to the randomness of diffusion models



■ Problem of existing works:

There are two characteristics of diffusion-based IR models:

- **Unique ground truths** of the task of image IR



■ Problem of existing works:

There are two characteristics of diffusion-based IR models:

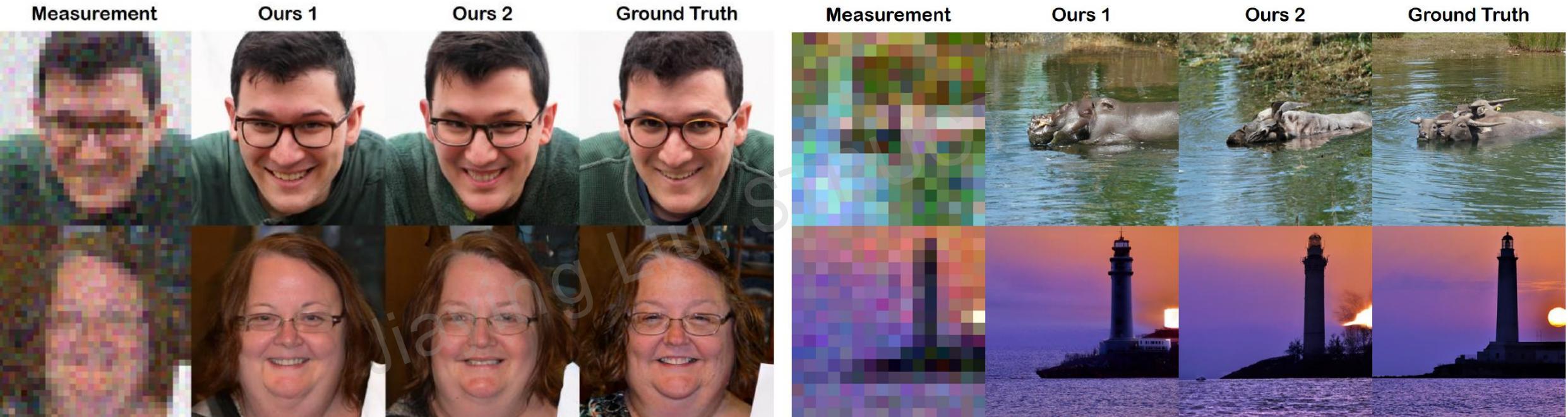
- Sampling **instability** due to the randomness of diffusion models

Vs. Unique ground truths of the task of image IR

It leads to the **problem**:

- The quality of each sample **cannot be guaranteed**

■ Problem:



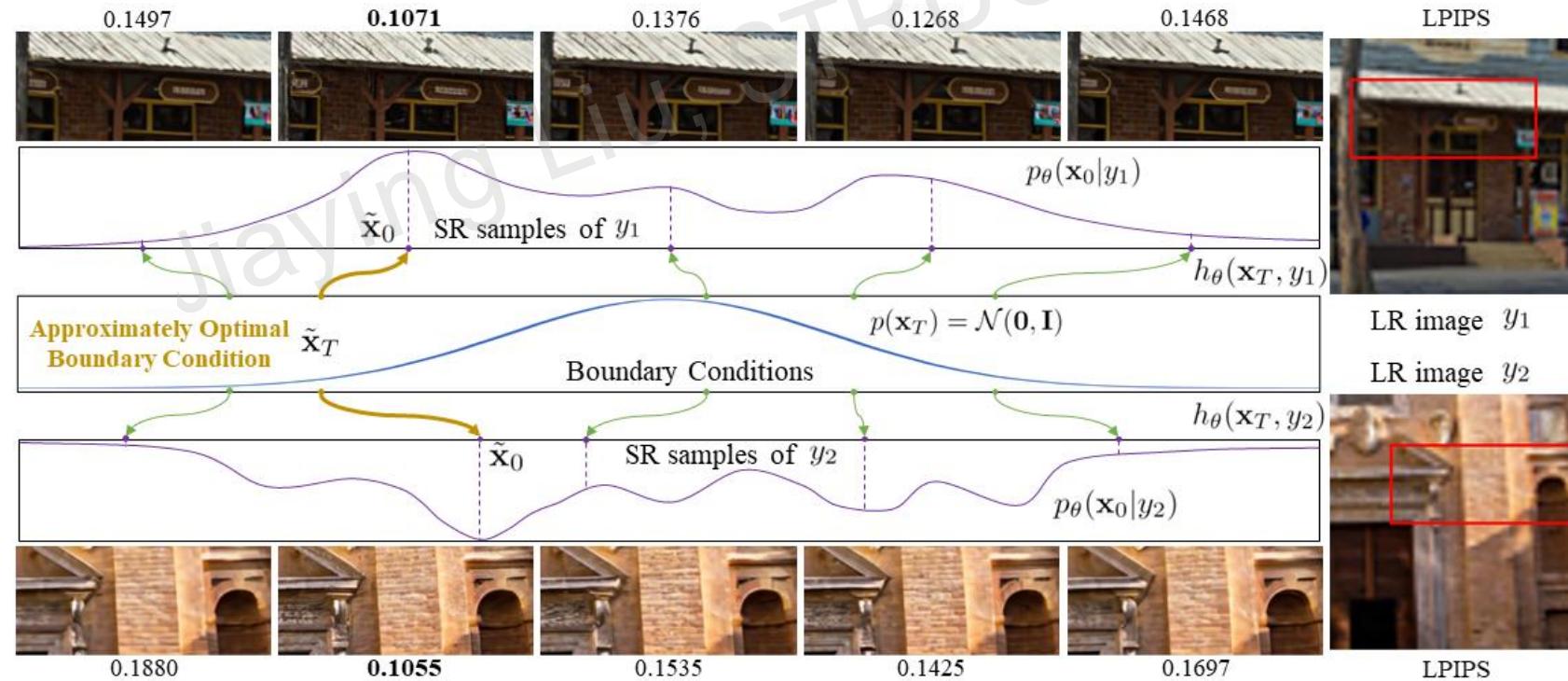
Diverse and reasonable, but there is **only one GT!**

Which one is **better**? How to **ensure the quality of the next sample**?

■ Our solution:

Solving Diffusion ODEs with Optimal Boundary Conditions for Better Image Super-Resolution

Yiyang Ma, Huan Yang, Wenhan Yang, Jianlong Fu, Jiaying Liu ICLR 2024



■ Existing Works:

Does not manage to solve the problem of instability.

■ Our Method:

A **stable sampling method** of diffusion-based SR models, **boosting all the diffusion-based SR models**.

- Analysis on the Approximately Optimal BC of Diffusion ODEs
- Obtaining the $\tilde{\mathbf{x}}_T$ with a refence set of image pairs
- Applying the $\tilde{\mathbf{x}}_T$ to all the LRs, achieving better results stably

■ Preliminary: Diffusion ODE sampler

The results are determined by the boundary condition \mathbf{x}_T .

(i.e., the Gaussian noise at the start of the sampling process).

e.g., DDIM [3], DPM-Solver [11]

Thus, the SR images can be a function of boundary condition (BC) and LR image:

$$\mathbf{x}_0 = h_\theta(\mathbf{x}_T, \mathbf{y})$$

[9] Jiaming Song, et al. Denoising Diffusion Implicit Models, ICLR, 2021

[10] Cheng Lu, et al. DPM-Solver: A Fast ODE Solver for Diffusion Probabilistic Model Sampling in Around 10 Steps, 2022

■ Our Solution: Optimal Boundary Condition \mathbf{x}_T^*

The BC that can generate the SR image which is the most close to the HR image:

$$\mathbf{x}_0^* = h_\theta(\mathbf{x}_T^*, \mathbf{y}) , \text{ where } \mathbf{x}_0^* = \arg \max_{\mathbf{x}_0} p_\theta(\mathbf{x}_0 | \mathbf{y})$$

We prove that the \mathbf{x}_T^* is approximately consistent to different LRs:

$$\mathbf{x}_T^* = \arg \max_{\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} p_\theta(h_\theta(\mathbf{x}_T, \mathbf{y})) \approx \arg \max_{\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} p_\theta(h_\theta(\mathbf{x}_T, \mathbf{y}_i)), \forall \mathbf{y}_i \in \mathcal{C}$$

■ Find an Approximately Optimal BC $\tilde{\mathbf{x}}_T$

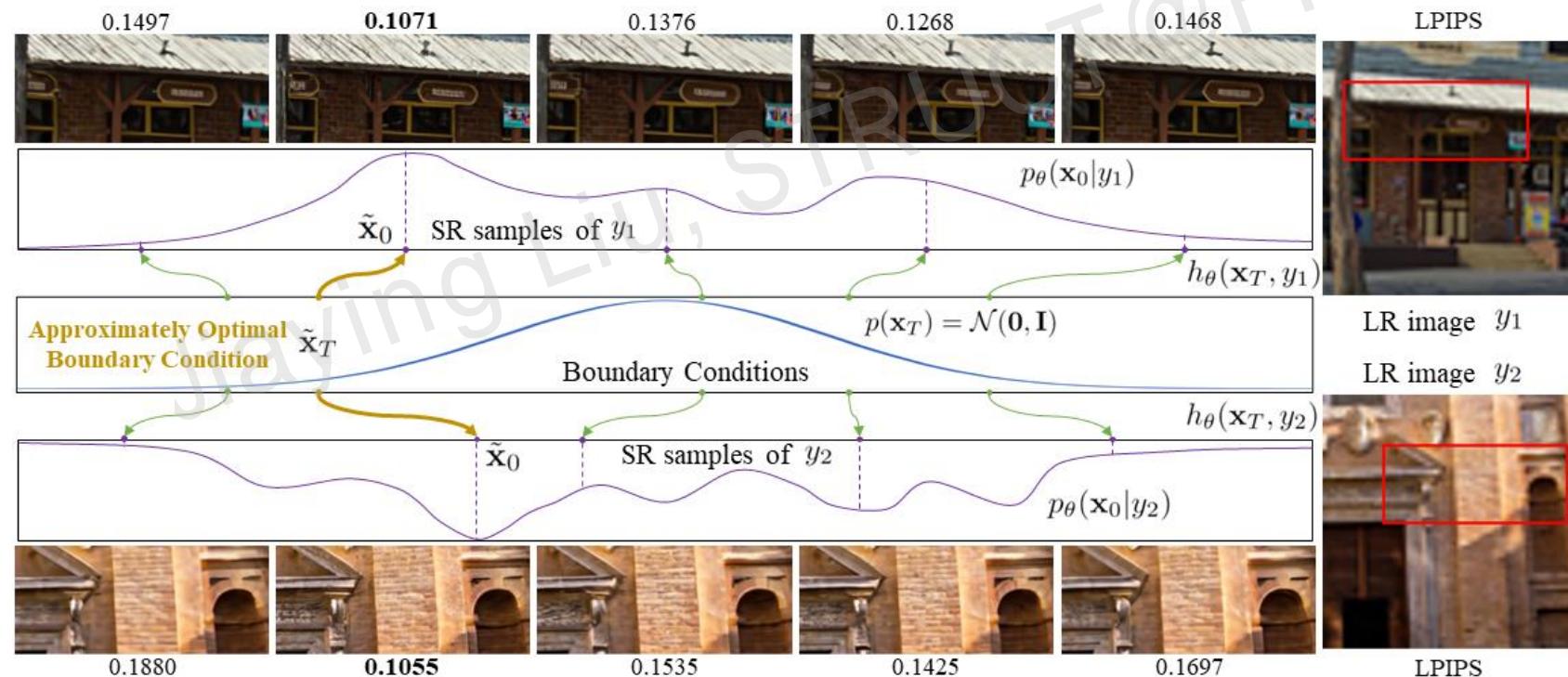
We employ a perceptual distance between SR and HR as an approximated implementation of the likelihood.

We build a reference set and calculate the $\tilde{\mathbf{x}}_T$ on it:

$$\tilde{\mathbf{x}}_T \approx \arg \min_{\mathbf{x}_T \sim \mathcal{N}(\mathbf{0}, \mathbf{I})} \sum_{i=1}^R M(h_\theta(\mathbf{x}_T, \mathbf{y}_i), \mathbf{z}_i)$$

■ Characteristic of Approximately Optimal BC $\tilde{\mathbf{x}}_T$

As we have discussed before, the $\tilde{\mathbf{x}}_T$ is approximately shared:



The one which can be used to sample the best SR image.

■ Characteristic of Approximately Optimal BC $\tilde{\mathbf{x}}_T$

As we have discussed before, the $\tilde{\mathbf{x}}_T$ is approximately shared:



The one which can be used to sample the best SR image.

■ Experiments

- The proposed method can be implemented on different SR models.

The results of SR3 on bicubic-SR.

| Model (& sampling method) | DIV2k-test | | Urban100 | | BSD100 | |
|---------------------------|--------------------------|--------|----------|--------|---------|--------|
| | LPIPS ↓ | PSNR ↑ | LPIPS ↓ | PSNR ↑ | LPIPS ↓ | PSNR ↑ |
| ESRGAN | 0.1082 | 28.18 | 0.1226 | 23.04 | 0.1579 | 23.65 |
| RankSRGAN | 0.1171 | 27.98 | 0.1403 | 23.16 | 0.1714 | 23.80 |
| SRDiff | 0.1286 | 28.96 | 0.1391 | 23.88 | 0.2046 | 24.17 |
| SR3 | DDPM-1000 | 0.1075 | 28.75 | 0.1165 | 24.33 | 0.1555 |
| | DDPM-250 | 0.1142 | 28.95 | 0.1181 | 24.41 | 0.1621 |
| | DDPM-100 | 0.1257 | 29.16 | 0.1232 | 24.51 | 0.1703 |
| | DPMS-20 | 0.1653 | 27.25 | 0.1413 | 23.46 | 0.2037 |
| | DDIM-50 | 0.1483 | 28.55 | 0.1333 | 24.16 | 0.1823 |
| | DDIM-100 | 0.1571 | 28.16 | 0.1335 | 24.05 | 0.1950 |
| | DPMS-20 + \tilde{x}_T | 0.1210 | 27.45 | 0.1179 | 23.57 | 0.1687 |
| | DDIM-50 + \tilde{x}_T | 0.1053 | 28.65 | 0.1164 | 24.26 | 0.1552 |
| | DDIM-100 + \tilde{x}_T | 0.1032 | 28.48 | 0.1136 | 24.12 | 0.1505 |

■ Experiments

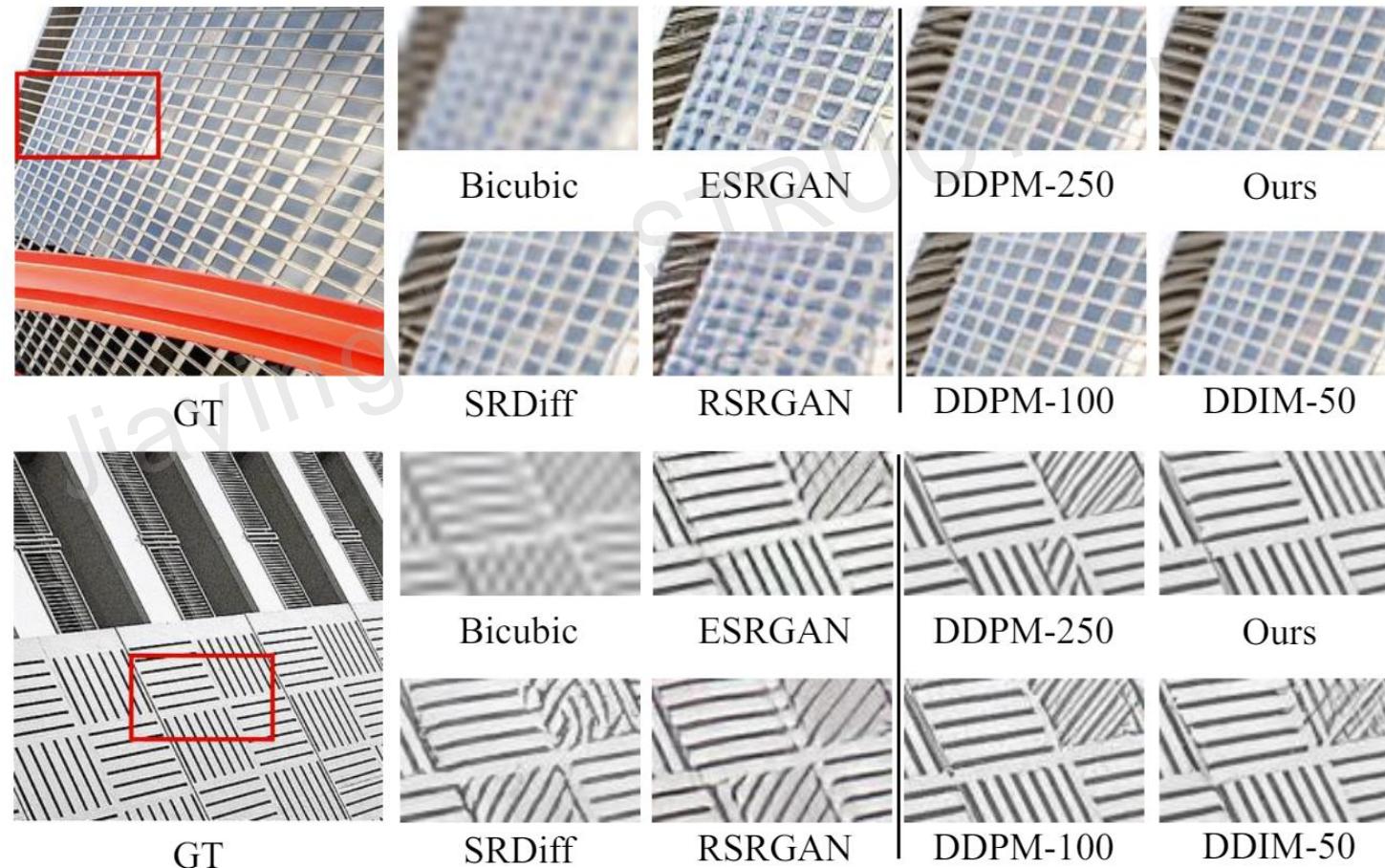
- The proposed method can be implemented on different SR models.

The results of StableSR and DiffIR on real world-SR.

| Model (& sampling method) | DIV2k-test | | | RealSR | | |
|---------------------------|-------------------------|---------------|---------------|---------------|---------------|---------------|
| | DISTS ↓ | LPIPS ↓ | PSNR ↑ | DISTS ↓ | LPIPS ↓ | PSNR ↑ |
| RealSR | 0.3051 | 0.5148 | 22.52 | 0.2532 | 0.3673 | 26.30 |
| BSRGAN | 0.2253 | 0.3416 | 22.13 | 0.2057 | 0.2582 | 25.52 |
| DASR | 0.2340 | 0.3444 | 22.02 | 0.2113 | 0.3014 | 26.32 |
| Real-ESRGAN | 0.2108 | 0.3109 | 22.36 | 0.2020 | 0.2511 | 25.12 |
| KDSR-GAN | 0.2022 | 0.2840 | 22.92 | 0.2006 | 0.2425 | 26.09 |
| StableSR | DDPM-200 | 0.2010 | 0.3189 | 19.42 | 0.2210 | 0.3065 |
| | DDIM-50 | 0.2217 | 0.3629 | 18.82 | 0.2336 | 0.3536 |
| | DDIM-50 + \tilde{x}_T | 0.2046 | 0.3169 | 19.55 | 0.2164 | 0.2999 |
| DiffIR | D-4 | 0.1773 | 0.2360 | 22.94 | 0.2076 | 0.2604 |
| | D-4 + \tilde{x}_T | 0.1772 | 0.2357 | 22.95 | 0.1993 | 0.2419 |

■ Experiments

■ Subjective results comparisons:



■ Summary of the method

We propose a **stable sampling method** of diffusion-based SR models,
boosting all the diffusion-based SR models.

- We analyze the characteristics of BCs.
- We propose a method of approximating an optimal BC.
- The proposed method can be utilized in different diffusion-based SR models without any external training.



02

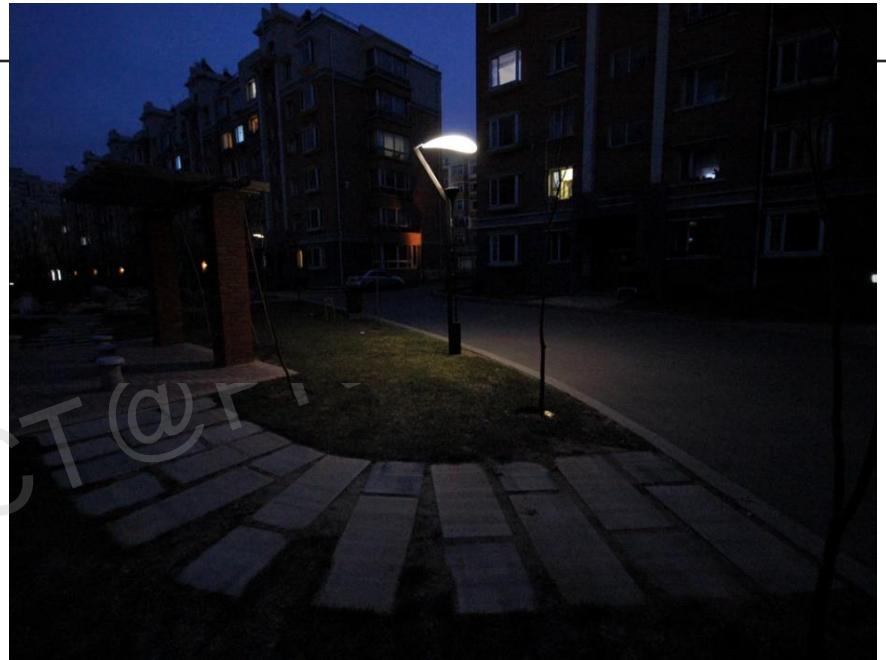
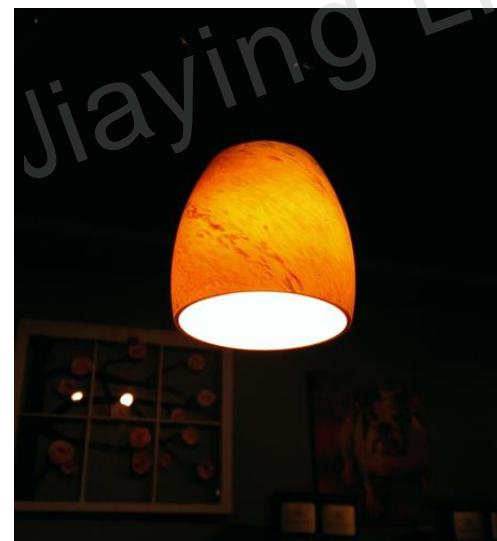
Diffusion for Low-light Enhancement

Diffusion for Low-Light Enhancement

53

■ Low visibility

- Details are buried due to degraded contrast and low illumination

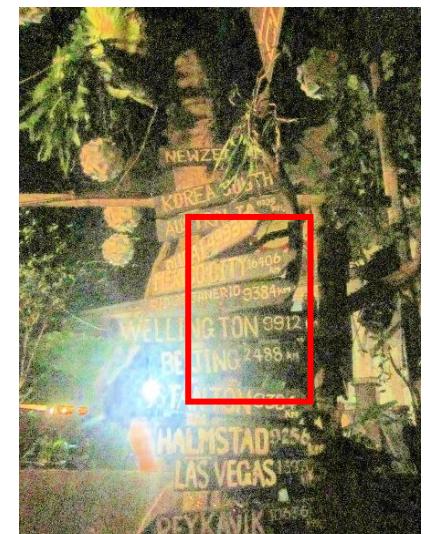
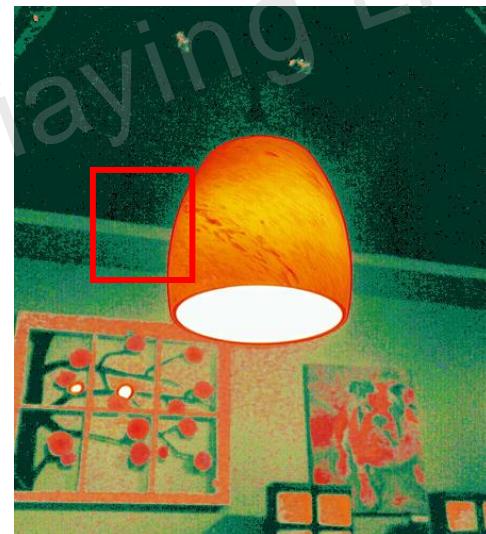


Diffusion for Low-Light Enhancement

54

■ Intensive noises

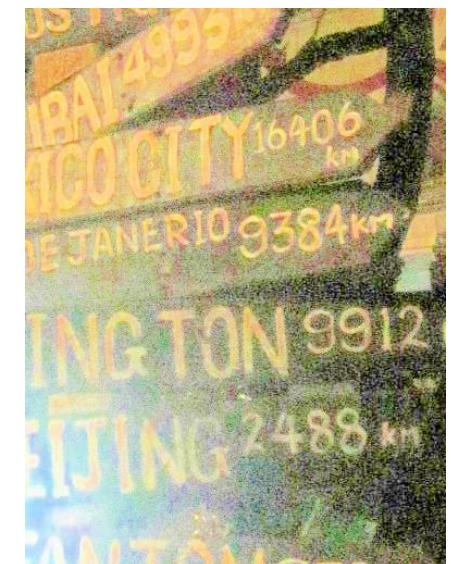
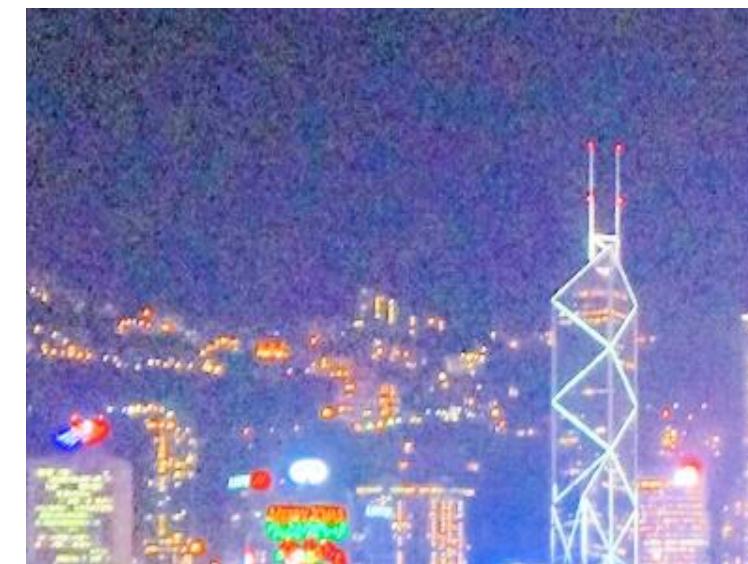
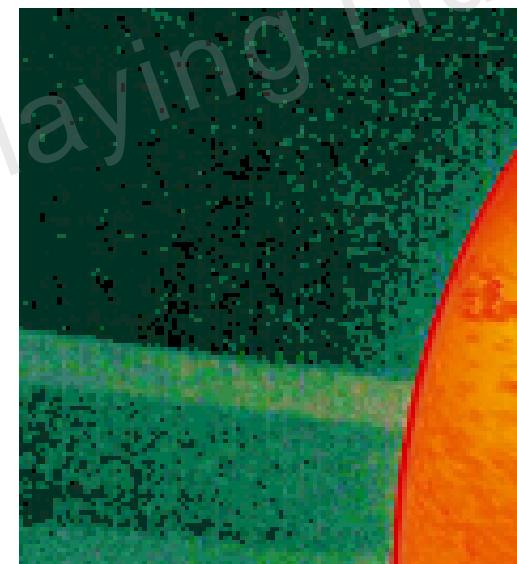
- After simple operations,
e.g. histogram equalization,
noises become visible



Diffusion for Low-Light Enhancement

55

- Intensive noises
 - After simple operations,
e.g. histogram equalization,
noises become visible

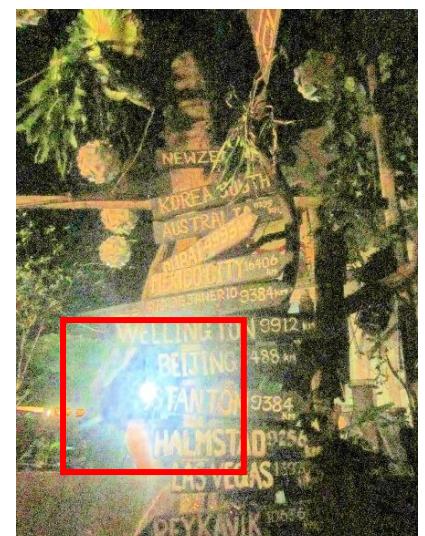


Diffusion for Low-Light Enhancement

56

■ Non-uniform illumination

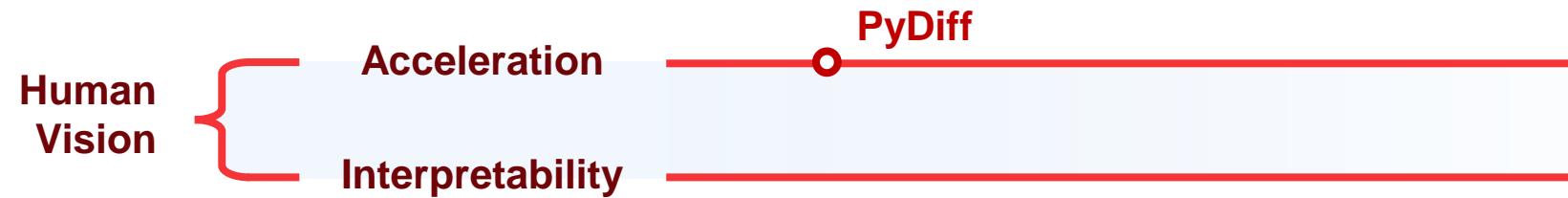
- Under-exposures
- Over-exposures



Diffusion for Low-Light Enhancement

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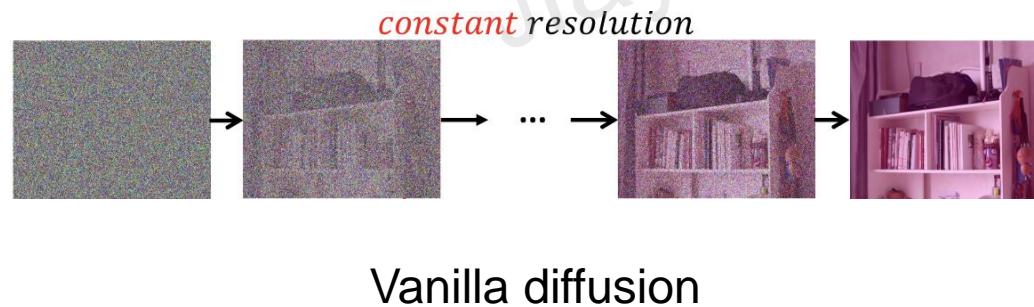


■ Pyramid Diffusion Model (PyDiff)

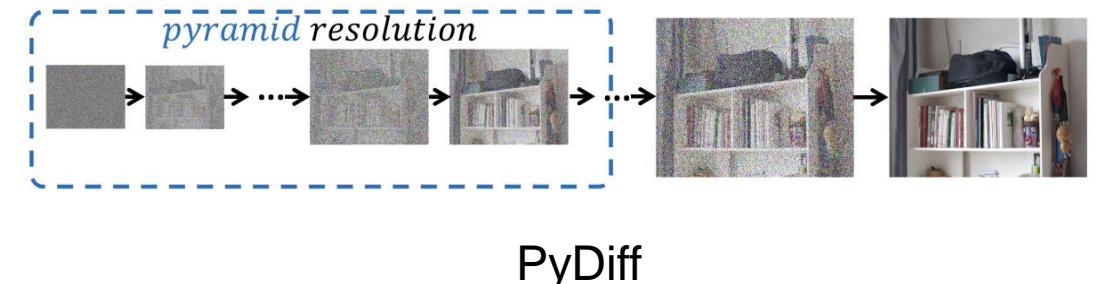
Problems of diffusion-based low-light enhancement

- Slow inference due to high resolution

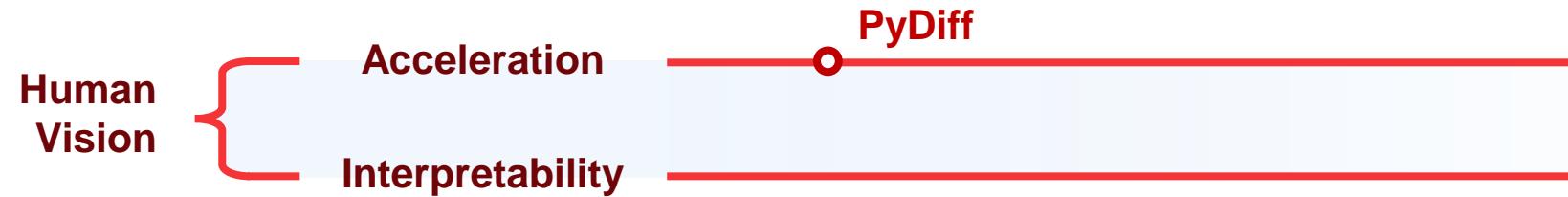
→ Sampling in a **pyramid** resolution style



Vanilla diffusion



PyDiff

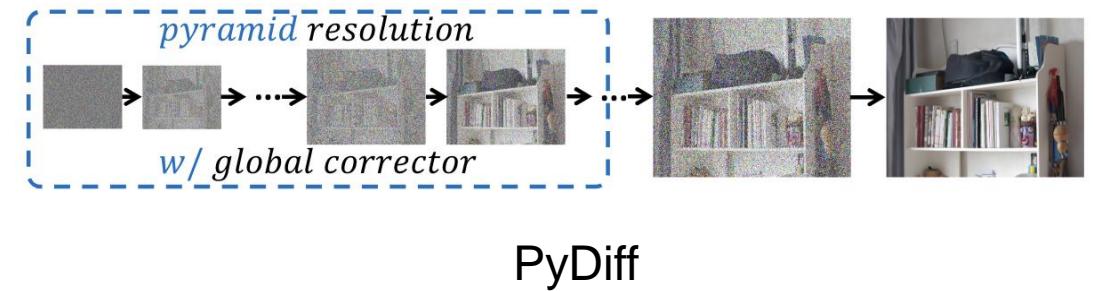
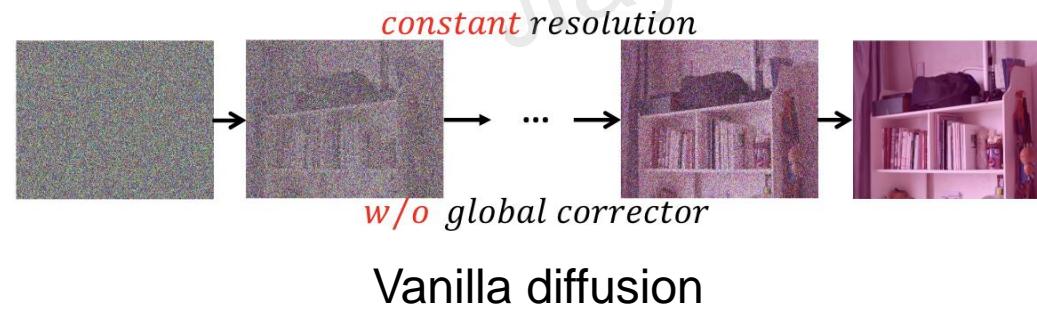


■ Pyramid Diffusion Model (PyDiff)

Problems of diffusion-based low-light enhancement

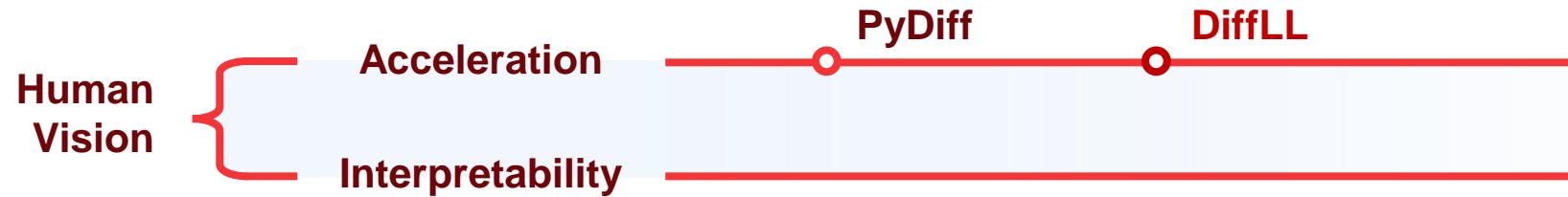
- Slow inference due to high resolution
- Global degradation (e.g., RGB shift)

- Sampling in a **pyramid resolution style**
- **A global corrector**



Diffusion for Low-Light Enhancement

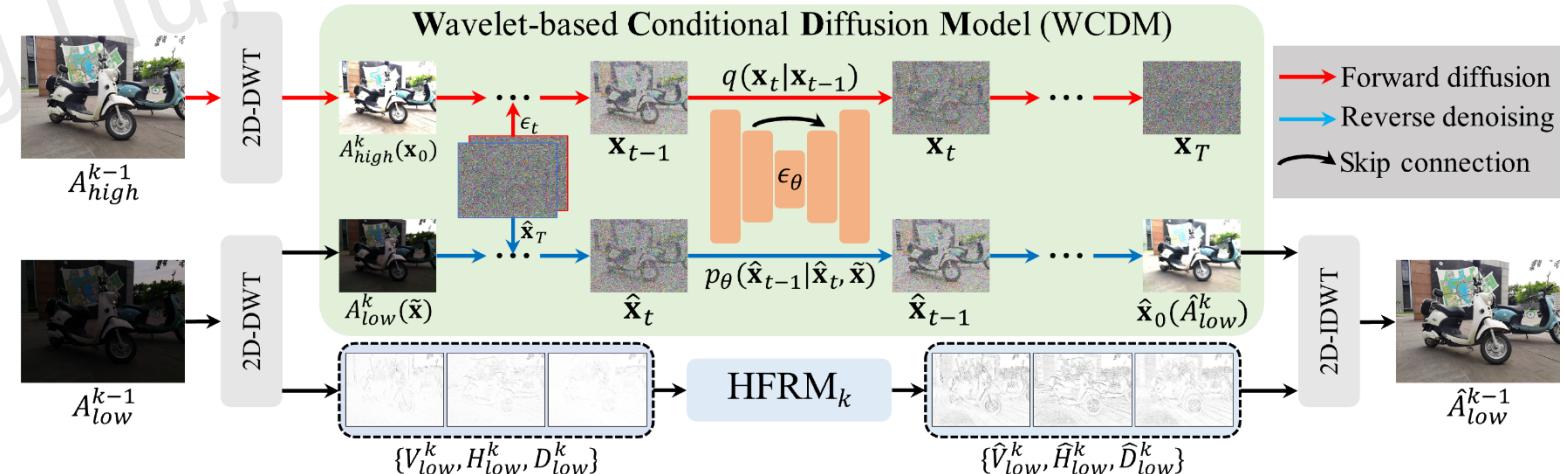
60



■ DiffLL

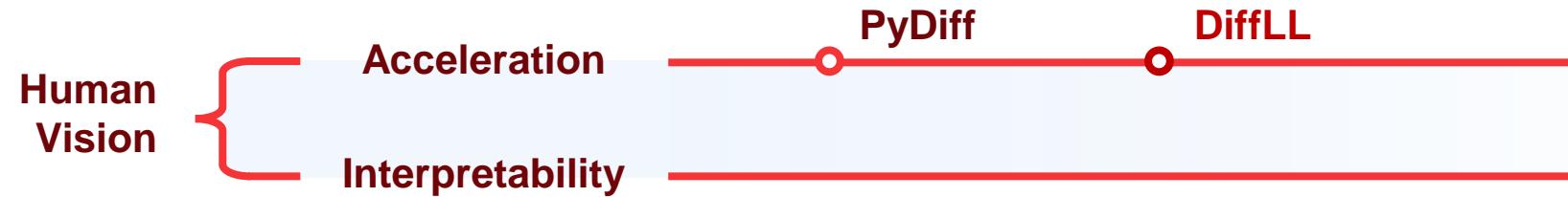
Wavelet transformation: accelerate inference

- *Diffusion-based WCDM*
restores **the average coefficient** A_{low}^k
- *Conv-based HFRM*
restores **high-frequency coefficients** $\{V_{low}^k, H_{low}^k, D_{low}^k\}$

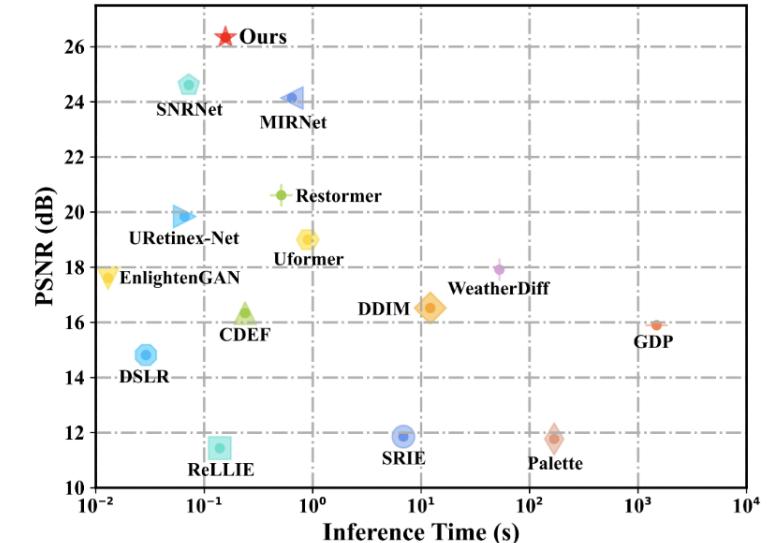


Diffusion for Low-Light Enhancement

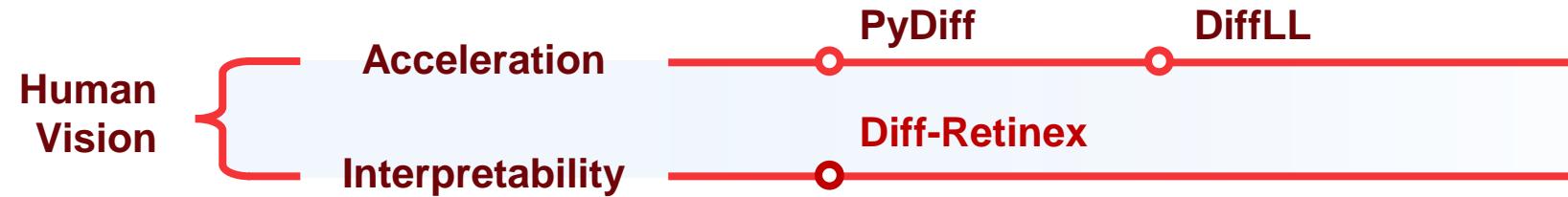
61



DiffLL

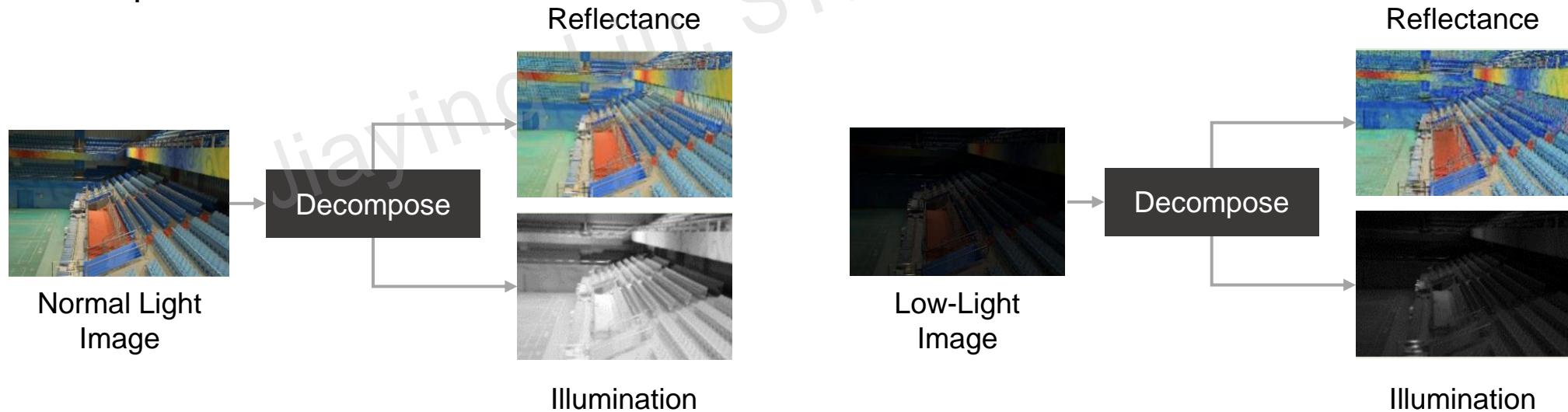


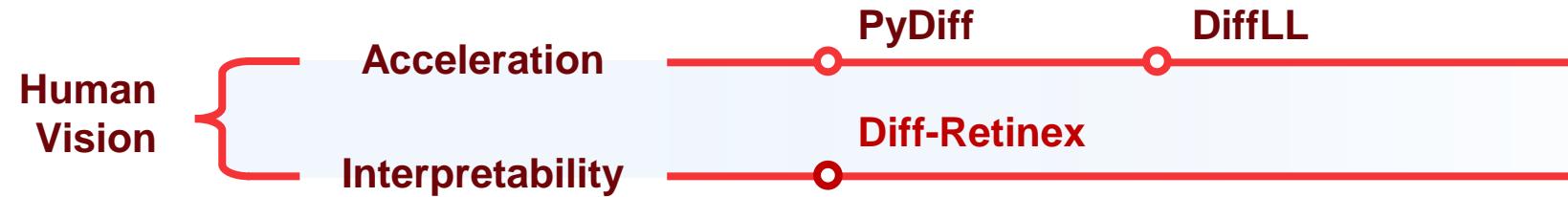
(a) Performance vs. Efficiency



■ Diff-Retinex

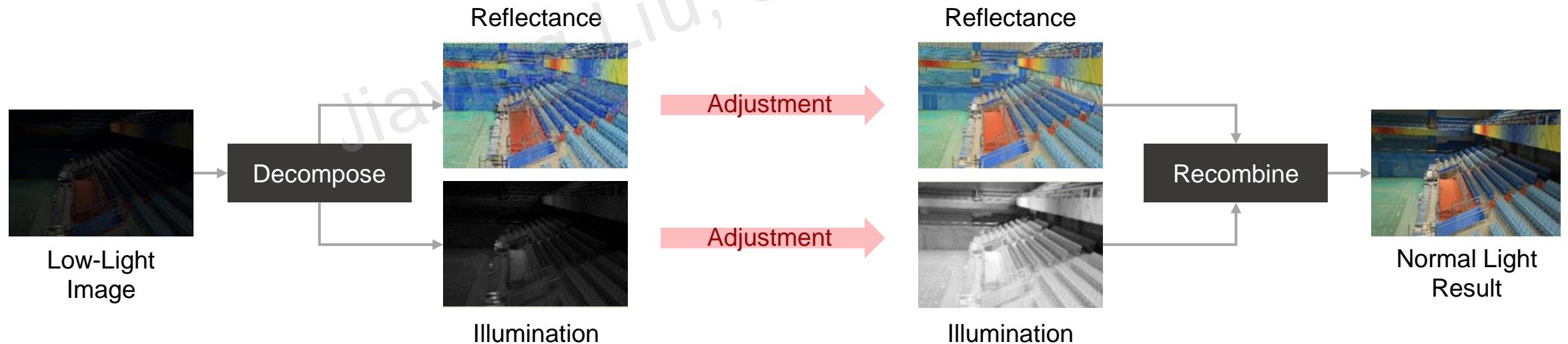
Retinex decomposition

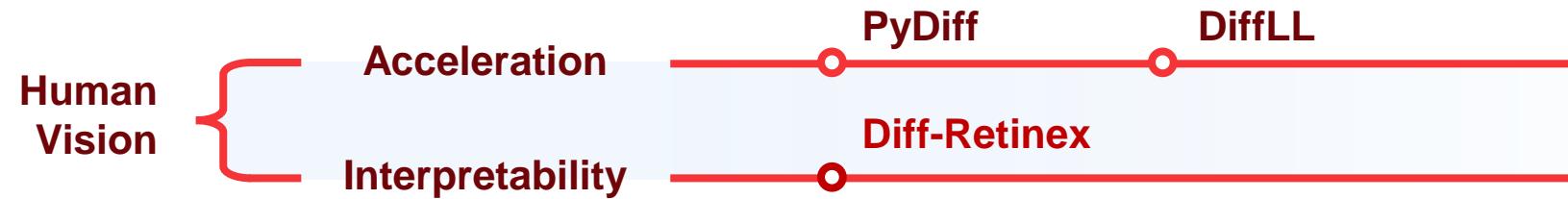




■ Diff-Retinex

Retinex-based low-light enhancement

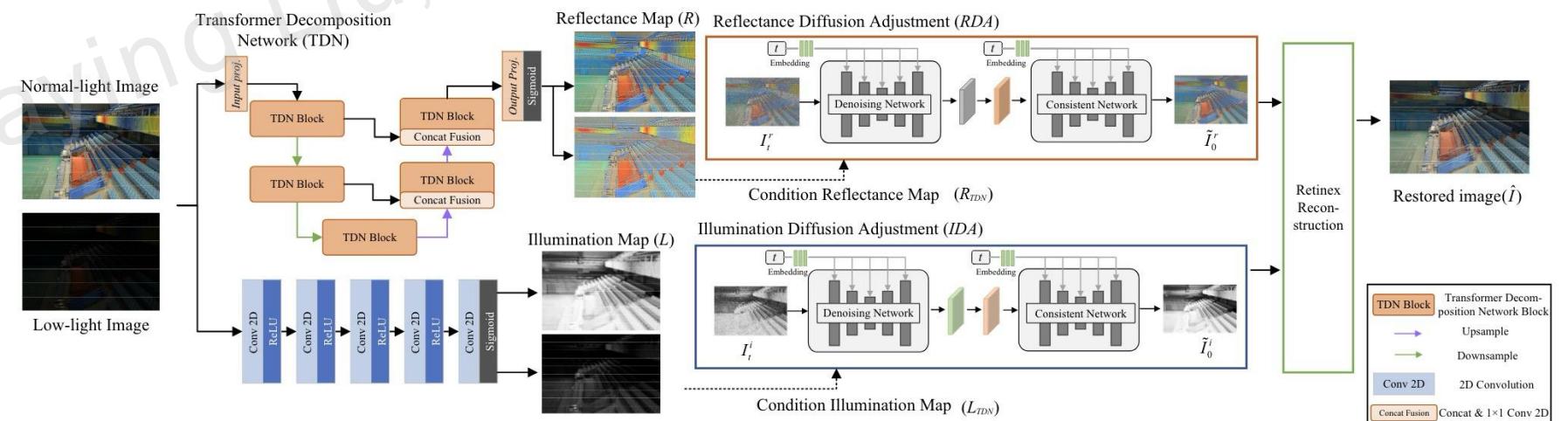




Diff-Retinex

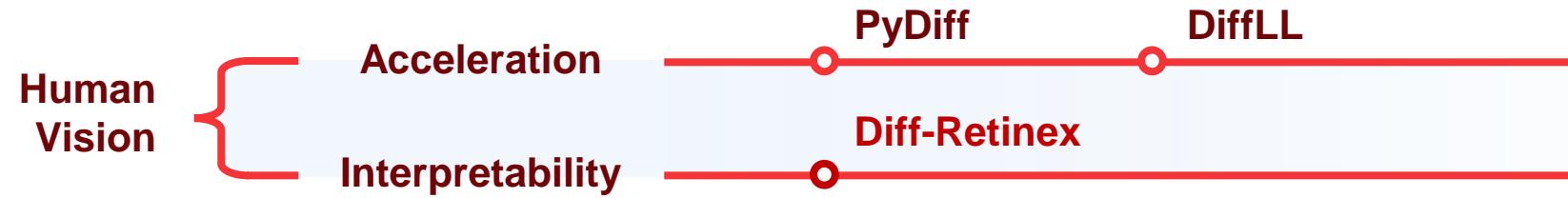
Physically explainable (**Retinex decomposition**) + Diffusion

- Diffusion for adjusting reflectance and illumination

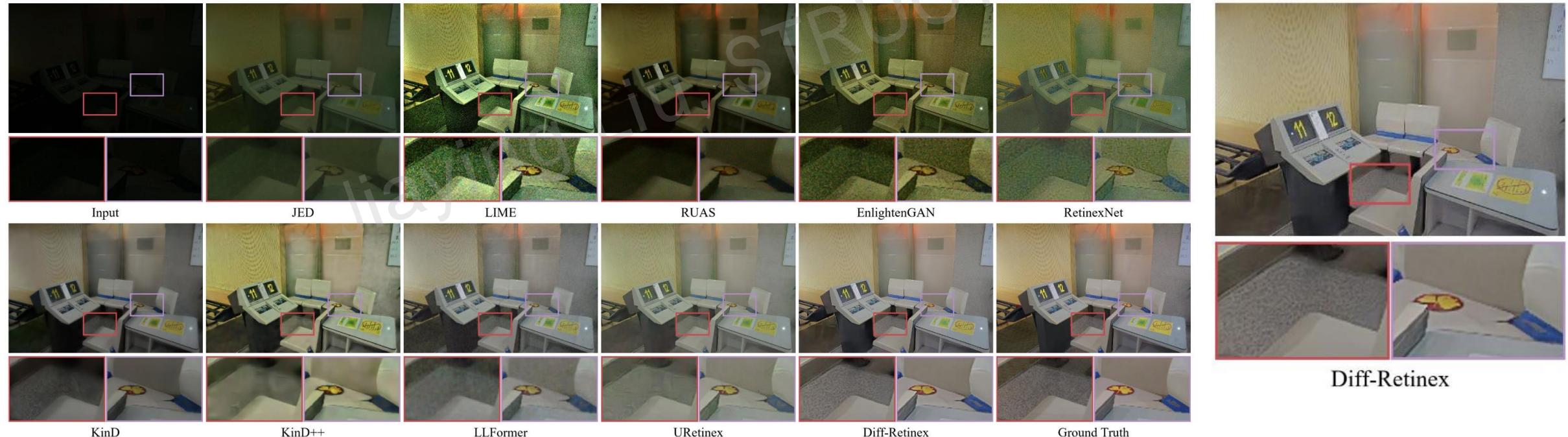


Diffusion for Low-Light Enhancement

65

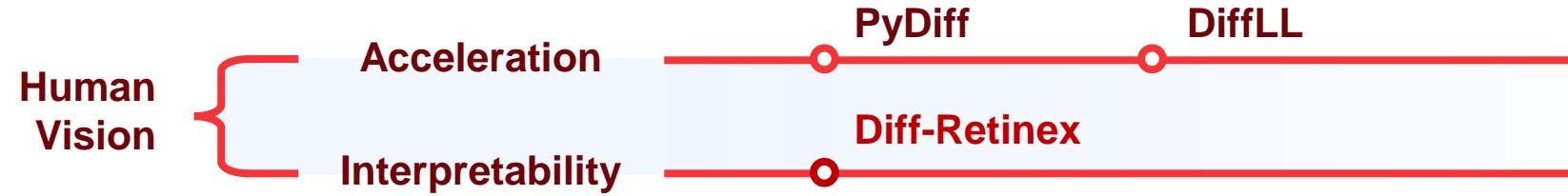


Diff-Retinex



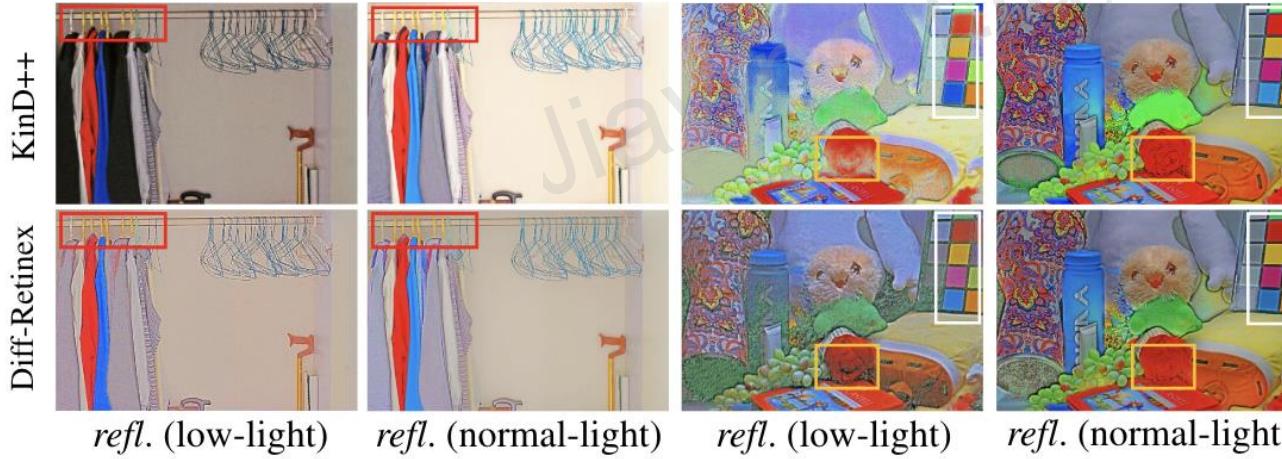
Diffusion for Low-Light Enhancement

66



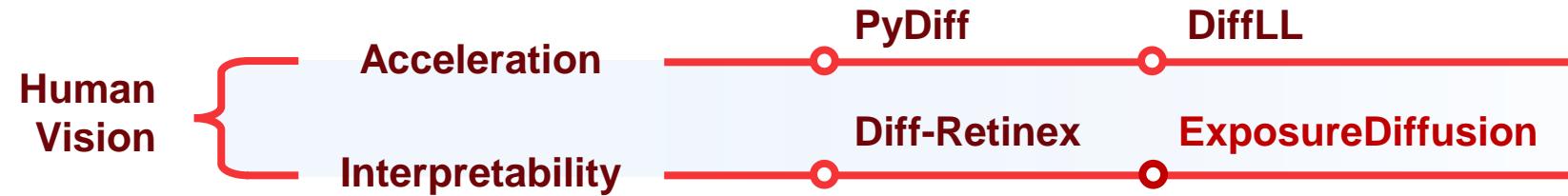
Diff-Retinex

Results of Retinex-decomposition



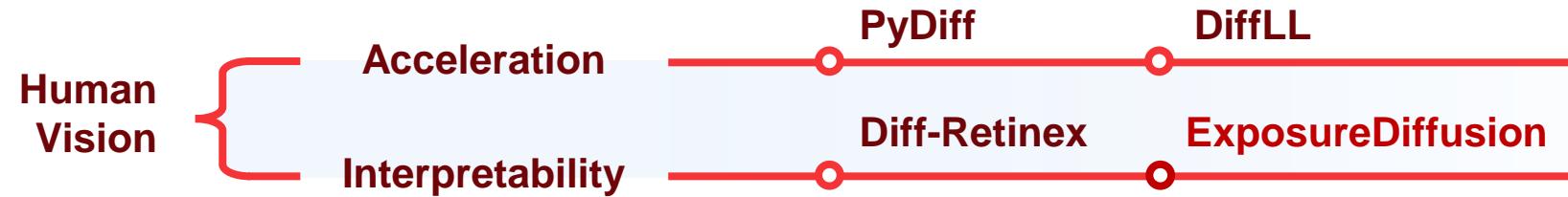
Visualization of restoring the reflectance map





■ ExposureDiffusion

Physically explainable (**exposure process**) + Diffusion



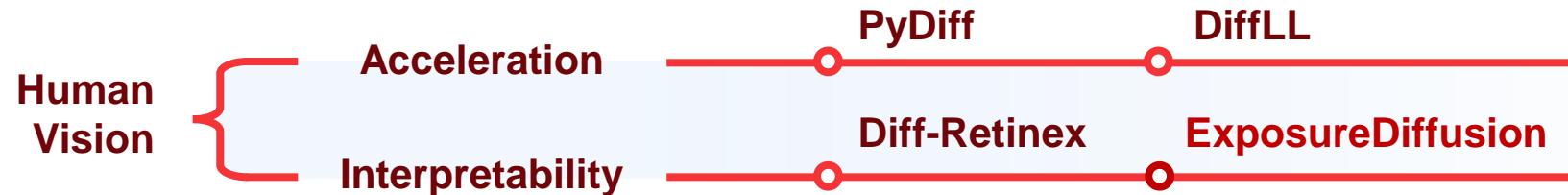
■ ExposureDiffusion

Physically explainable (**exposure process**) + Diffusion

- A raw image can be formulated as

$$X_t = \lambda_t K I + K N_p + N_{ind}$$

- λ_t exposure time $(\lambda_t I + N_p) \sim \mathcal{P}(\lambda_t I)$
- K overall system gain **Poisson distribution**
- I rate of the photoelectrons
- N_p photon shot noise
- N_{ind} signal independent noise



■ ExposureDiffusion

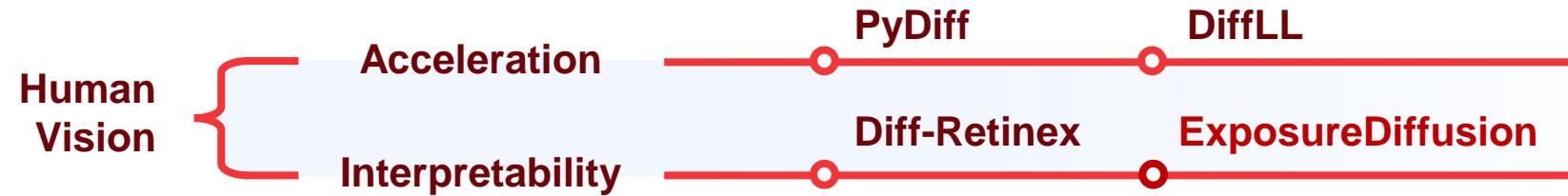
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- N_p photon shot noise
- N_{ind} signal independent noise

| | Conditional diffusion | Exposure diffusion |
|---------------------|------------------------------------------------------------------------------------------------|-----------------------------------------------------------------------------------------------------------------------------------------|
| Objective | maximizing $p_\Theta(X Y)$ | minimizing KL divergence with the real exposure process |
| Initial state X_T | $X_T \sim \mathcal{N}(0, 1)$ | $X_T \sim q(X_T)$ |
| Assumption | $q(X_t X_{t-1}) := \mathcal{N}(X_t; \sqrt{1 - \beta_t} X_{t-1}, \beta_t \mathbf{I})$ | $q(X_{t-1} X_t, X_{ref}) := \mathcal{P}\left(\frac{X_{t-1}-X_t}{K}; \frac{(\lambda_{t-1}-\lambda_t)X_{ref}}{\lambda_{ref}K}\right)$ |
| Reverse process | $p_\Theta(X_{t-1} X_t, Y) := \mathcal{N}(X_{t-1}; \mu_\Theta(X_t, Y, t), \sigma^2 \mathbf{I})$ | $p_\Theta(X_{t-1} X_t) := \mathcal{P}\left(\frac{X_{t-1}-X_t}{K}; \frac{(\lambda_{t-1}-\lambda_t)F_\Theta(X_t)}{\lambda_{ref}K}\right)$ |
| Training | The expectation over $q(X_t X_0)$ | The expectation over $p_\Theta(X_t X_T)$ |

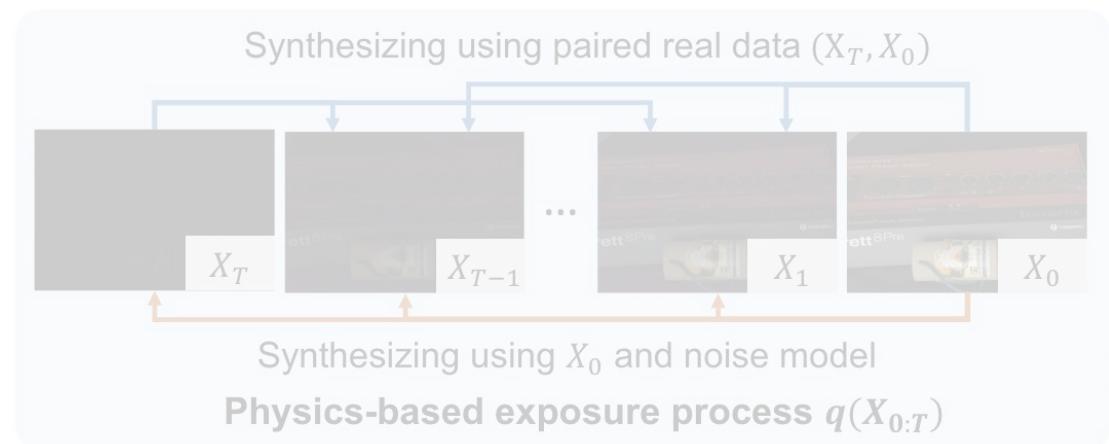
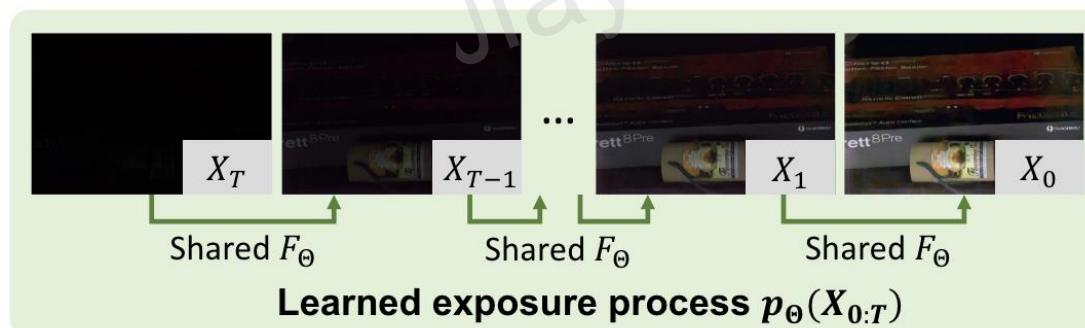


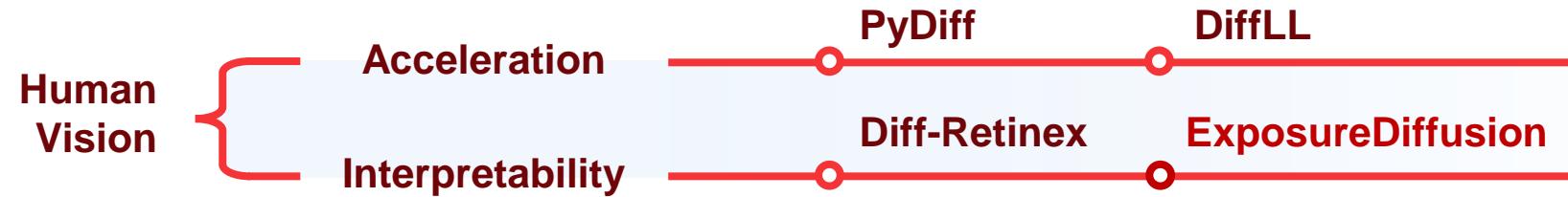
■ ExposureDiffusion

Physically explainable (**exposure process**) + Diffusion

- Simulate using a **shared** neural network
- In a progressive manner

Applicable to both **real-captured** and **synthetic** noise models

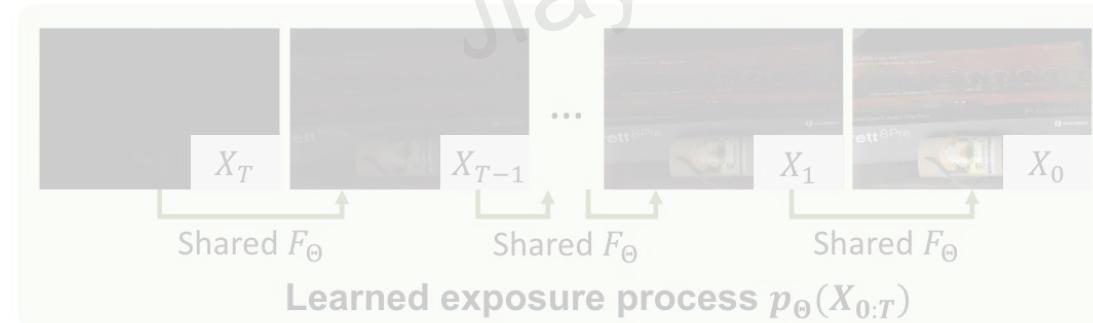




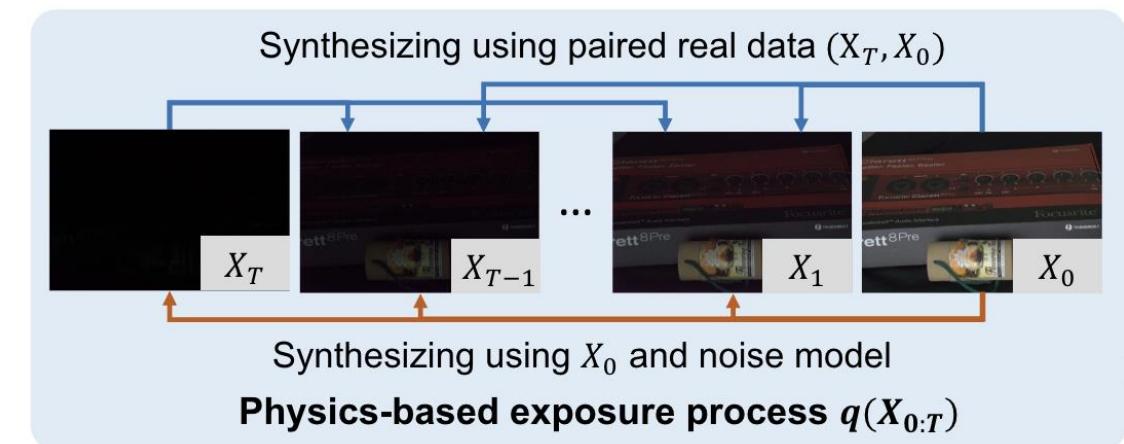
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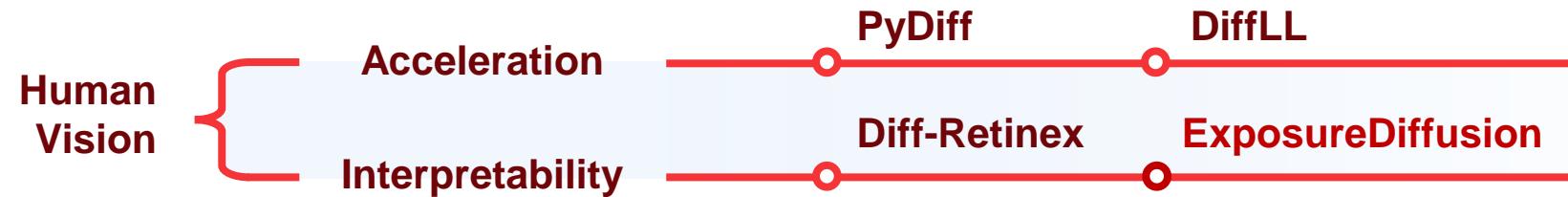


Applicable to both **real-captured** and **synthetic** noise models

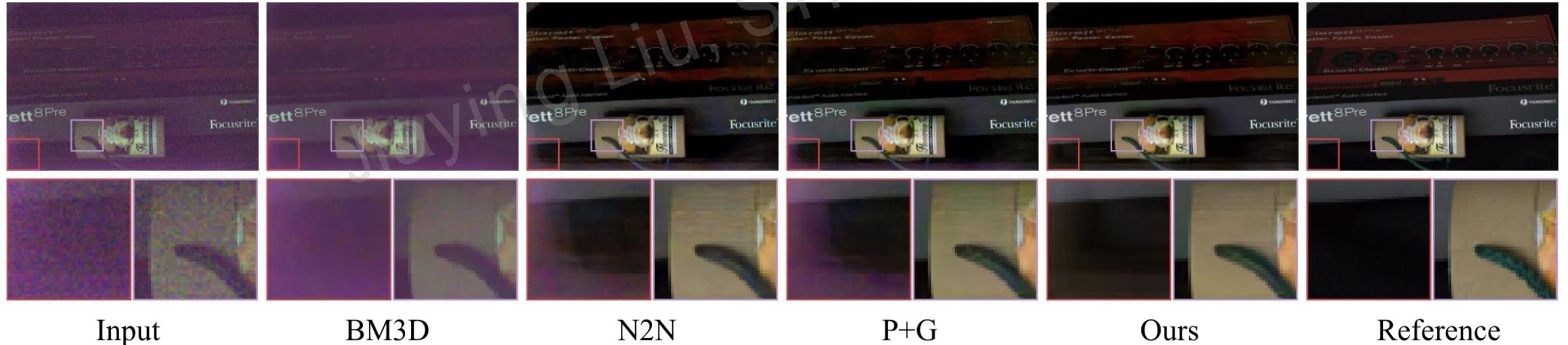


Diffusion for Low-Light Enhancement

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■ ExposureDiffusion



Diffusion for Low-Light Enhancement

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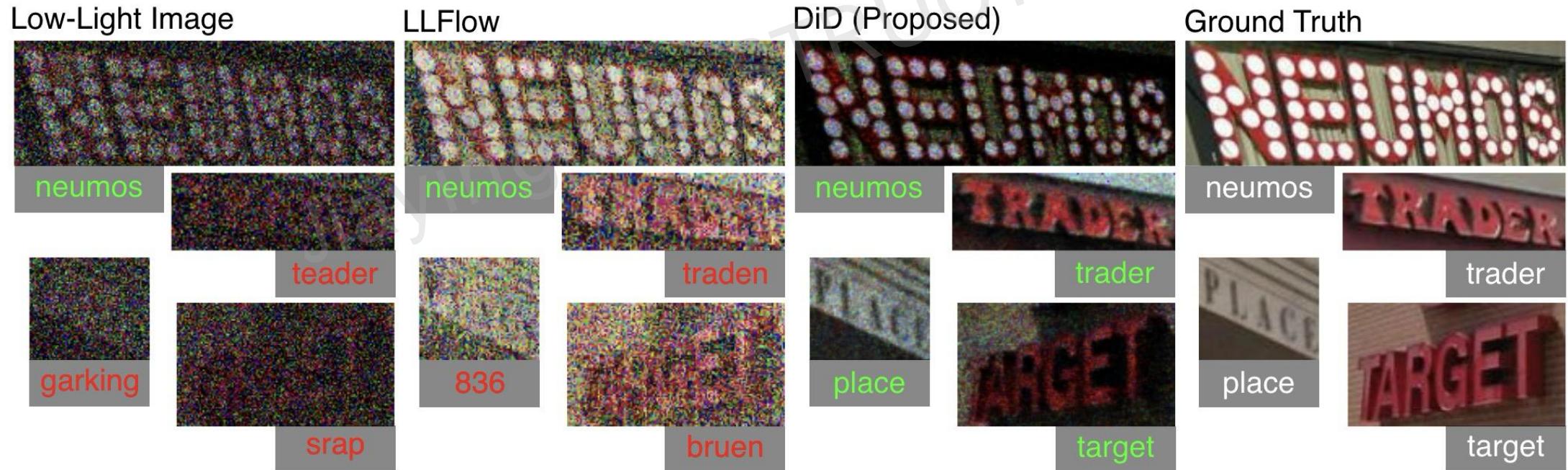


Diffusion for Low-Light Enhancement

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Diffusion in the Dark (DiD) for low-light text recognition

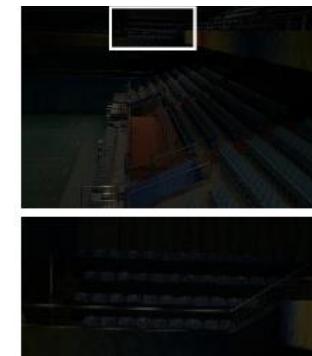


■ Problem of existing works

Cannot maintain a consistent brightening effect in dark and moderately dark images

Essential problem:

Lack a genuine
concept of illumination



Input

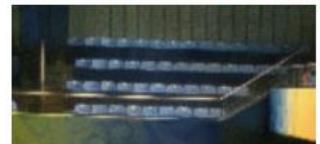


SCI-MIT^[1]



SCI-LOL^[1]

Over-exposure



Ours

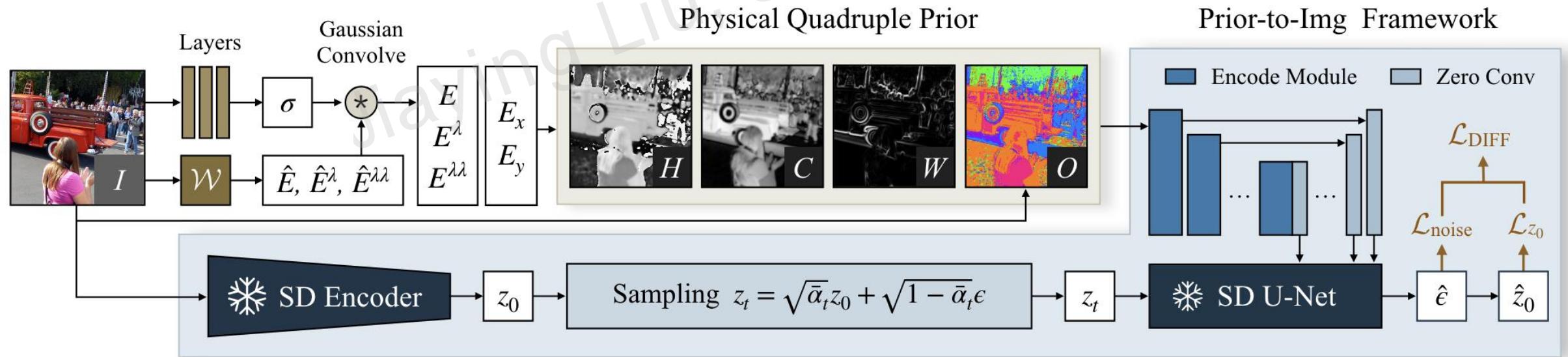
[1] Long Ma, et al. Toward fast, flexible, and robust low-light image enhancement. In CVPR, 2022.

■ Our solution:

Zero-Reference Low-Light Enhancement via Physical Quadruple Priors

Wenjing Wang, Huan Yang, Jianlong Fu, Jiaying Liu

CVPR 2024

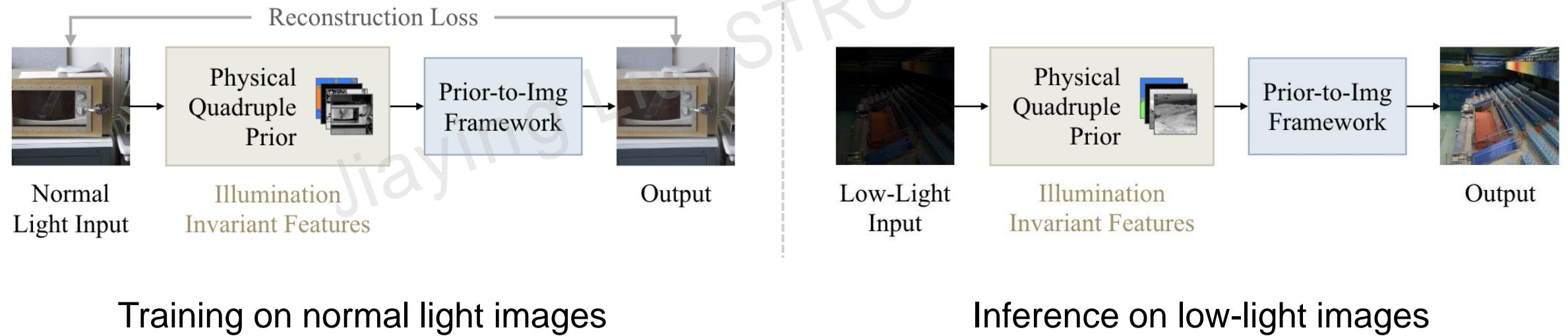


Diffusion for Low-Light Enhancement

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■ Our solution:

- Develop an illumination-invariant prior (**Physical Quadruple Prior**)
- Employ it as an intermediary between low-light and normal light images

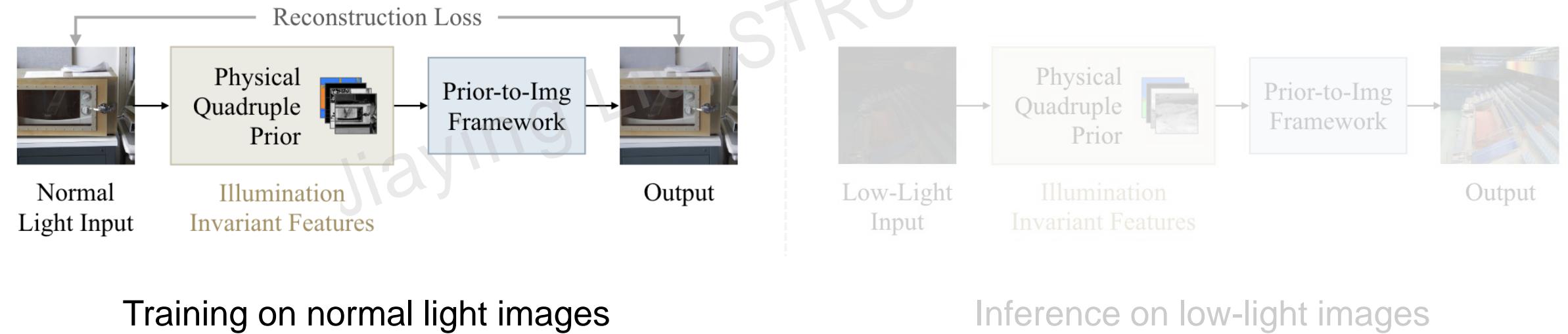


Diffusion for Low-Light Enhancement

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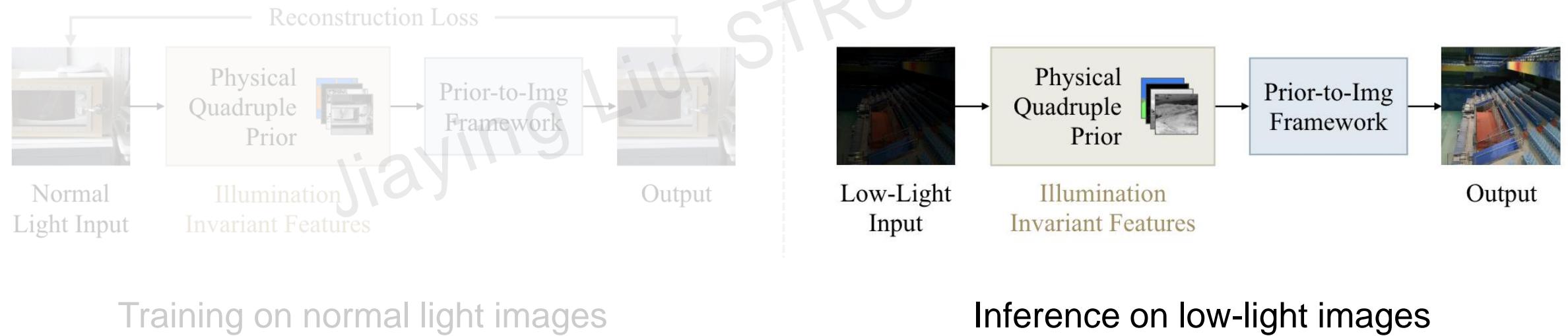
■ Our solution:

- Develop an illumination-invariant prior (**Physical Quadruple Prior**)
- Employ it as an intermediary between low-light and normal light images



■ Our solution:

- Develop an illumination-invariant prior (**Physical Quadruple Prior**)
- Employ it as an intermediary between low-light and normal light images



■ The Kubelka-Munk theory of light transfer

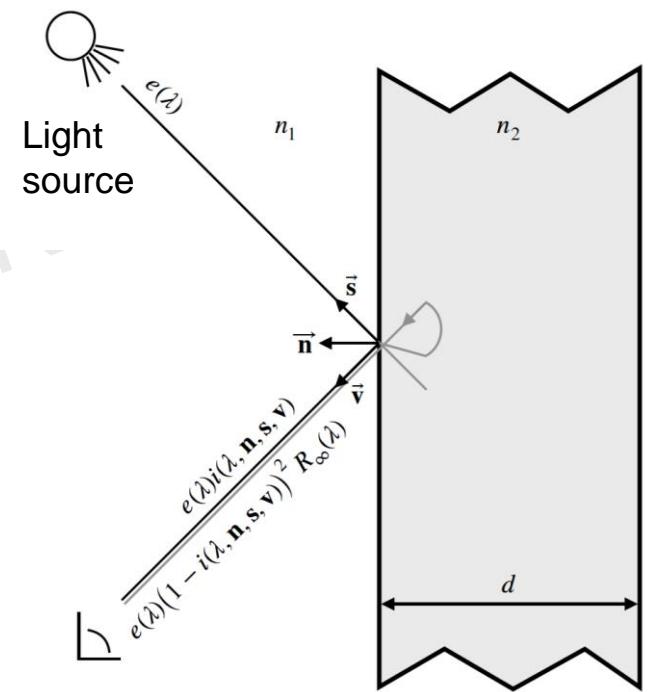
$$E(\lambda, \mathbf{x}) = e(\lambda, \mathbf{x}) ((1 - i(\mathbf{x}))^2 R_\infty(\lambda, \mathbf{x}) + i(\mathbf{x})),$$

$e(\lambda, \mathbf{x})$: the spectrum of the light source
 $i(\mathbf{x})$: the specular reflection
 $R_\infty(\lambda, \mathbf{x})$: the material reflectivity

} Illumination dependent
→ Illumination invariant

■ when the object is matte, i.e., $i(\mathbf{x}) \approx 0$

$$E(\lambda, \mathbf{x}) = e(\lambda, \mathbf{x}) R_\infty(\lambda, \mathbf{x}), \quad \text{the Retinex theory is a special case}$$



- Through simplifying assumptions, obtain a series of invariants
 - Assuming equal energy illumination

$$E(\lambda, \mathbf{x}) = \tilde{e}(\mathbf{x}) ((1 - i(\mathbf{x}))^2 R_\infty(\lambda, \mathbf{x}) + i(\mathbf{x}))$$

- Then we can obtain

$$\frac{E^\lambda}{E^{\lambda\lambda}} = \frac{\tilde{e}(\mathbf{x})(1 - i(\mathbf{x}))^2 R_\infty^\lambda}{\tilde{e}(\mathbf{x})(1 - i(\mathbf{x}))^2 R_\infty^{\lambda\lambda}} = \boxed{\frac{R_\infty^\lambda}{R_\infty^{\lambda\lambda}}}, \text{ Illumination Invariant}$$

- Set $H = \arctan(E^\lambda/E^{\lambda\lambda})$, we can find that H is illumination invariant

- Through simplifying assumptions, obtain a series of invariants
 - Assuming equal energy illumination

$$E(\lambda, \mathbf{x}) = \tilde{e}(\mathbf{x}) ((1 - i(\mathbf{x}))^2 R_\infty(\lambda, \mathbf{x}) + i(\mathbf{x})) \rightarrow H = \arctan(E^\lambda / E^{\lambda\lambda})$$

- Through simplifying assumptions, obtain a series of invariants
 - Assuming equal energy illumination

$$E(\lambda, \mathbf{x}) = \tilde{e}(\mathbf{x}) ((1 - i(\mathbf{x}))^2 R_\infty(\lambda, \mathbf{x}) + i(\mathbf{x})) \rightarrow H = \arctan(E^\lambda / E^{\lambda\lambda})$$

- Further assuming that the surface is matte

$$E(\lambda, \mathbf{x}) = \tilde{e}(\mathbf{x}) R_\infty(\lambda, \mathbf{x}), \rightarrow C = \log \left(\frac{(E^\lambda)^2 + (E^{\lambda\lambda})^2}{E(\lambda, \mathbf{x})^2} \right)$$

- Further assuming uniform illumination

$$E(\lambda, \mathbf{x}) = \bar{e} R_\infty(\lambda, \mathbf{x}), \rightarrow W = \tan \left(\left| \frac{\partial E(\lambda, \mathbf{x})}{\partial \mathbf{x}} \frac{1}{E(\lambda, \mathbf{x})} \right| \right)$$

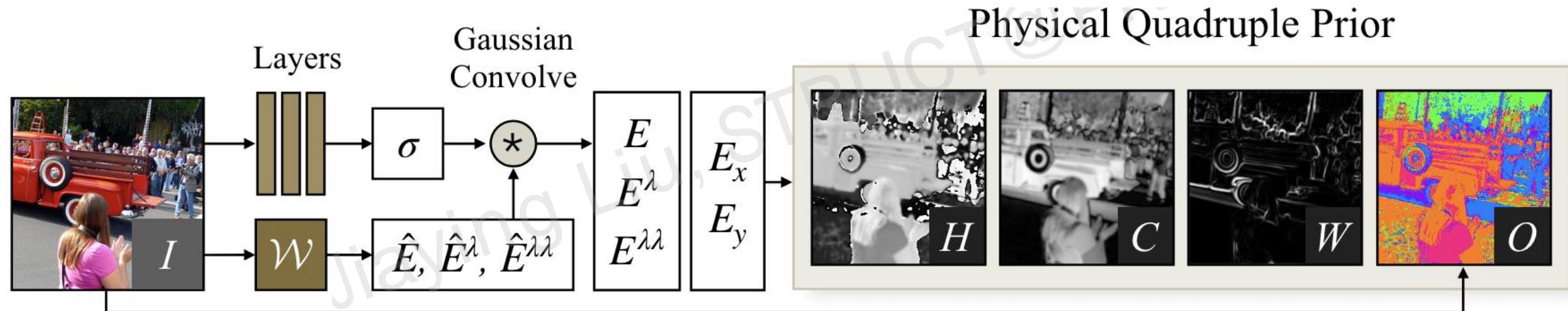
- Through simplifying assumptions, obtain a series of invariants

$$H = \arctan(E^\lambda / E^{\lambda\lambda}) \quad C = \log \left(\frac{(E^\lambda)^2 + (E^{\lambda\lambda})^2}{E(\lambda, \mathbf{x})^2} \right) \quad W = \tan \left(\left| \frac{\partial E(\lambda, \mathbf{x})}{\partial \mathbf{x}} \frac{1}{E(\lambda, \mathbf{x})} \right| \right)$$

- Assuming that illumination maintains the order of colors, we propose the order of the RGB three channels as a fundamental illumination-invariant feature
- $O(x, y) = [O_R(x, y), O_G(x, y), O_B(x, y)]$
- H, C, W , and O are concatenated in the channel dimension to form our **physical quadruple prior**

■ Learnable Illumination-Invariant Prior

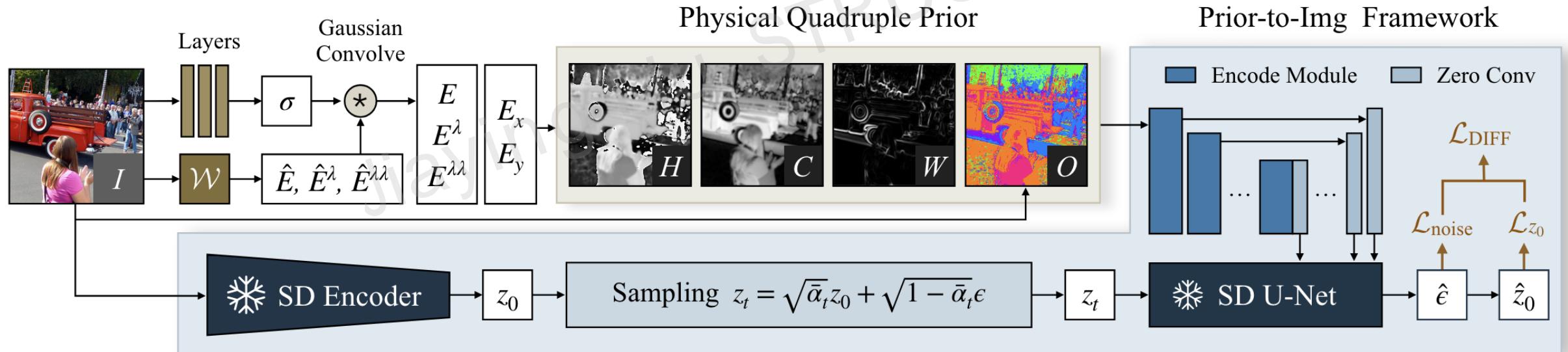
Learn within deep neural networks



■ Prior-to-Image via Diffusion Models

Leverage the knowledge of a pretrained image generative model

Use **Physical Quadruple Prior** as the condition of generation



■ Prior-to-Image via Diffusion Models

Issues of generative models: *slow convergence, detail degradation*

■ Prior-to-Image via Diffusion Models

Issues of generative models: *slow convergence, detail degradation*

■ Improvements to make it more suitable for image restoration task

1. Training objective

$$\mathcal{L}_{z_0} = \|z_0 - \hat{z}_0\|_2^2 = \|z_0 - \frac{z_t - \sqrt{1 - \bar{\alpha}_t}\hat{\epsilon}}{\sqrt{\bar{\alpha}_t}}\|_2^2$$

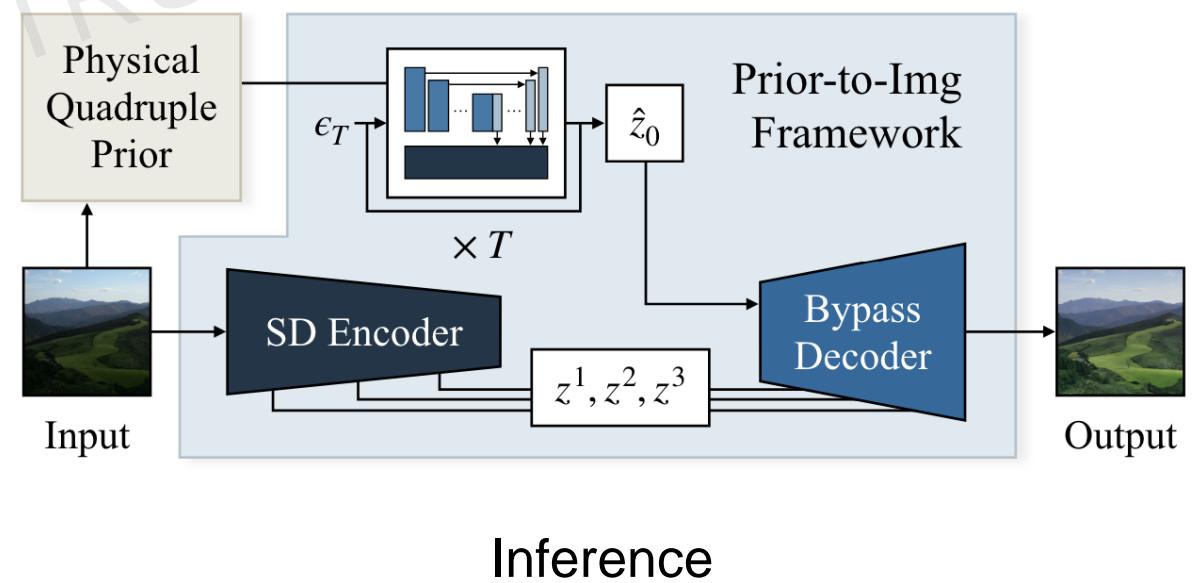
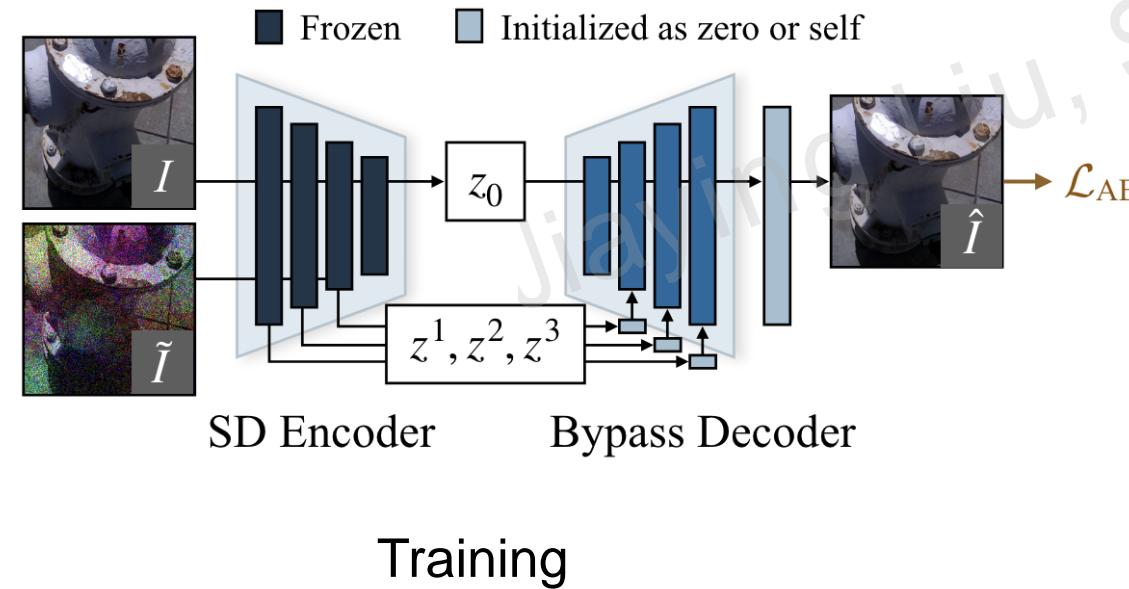
$$\mathcal{L}_{\text{noise}} = \|\epsilon - \hat{\epsilon}\|_2^2$$

$$\mathcal{L}_{\text{DIFF}} = \boxed{\mathcal{L}_{z_0}} + \boxed{\mathcal{L}_{\text{noise}}}$$

- Improvements to make it more suitable for image restoration task

2. Bypass decoder

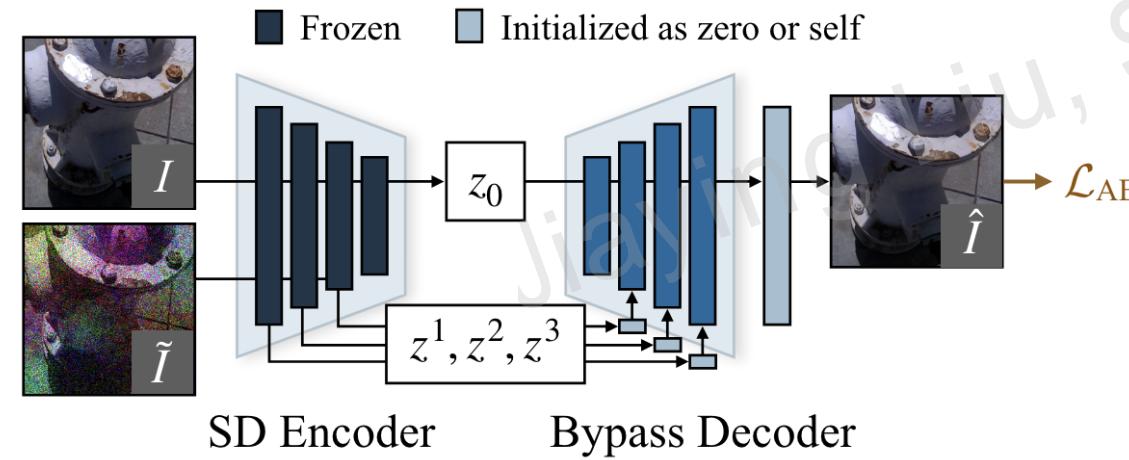
- Handle detail distortion through a new decoder



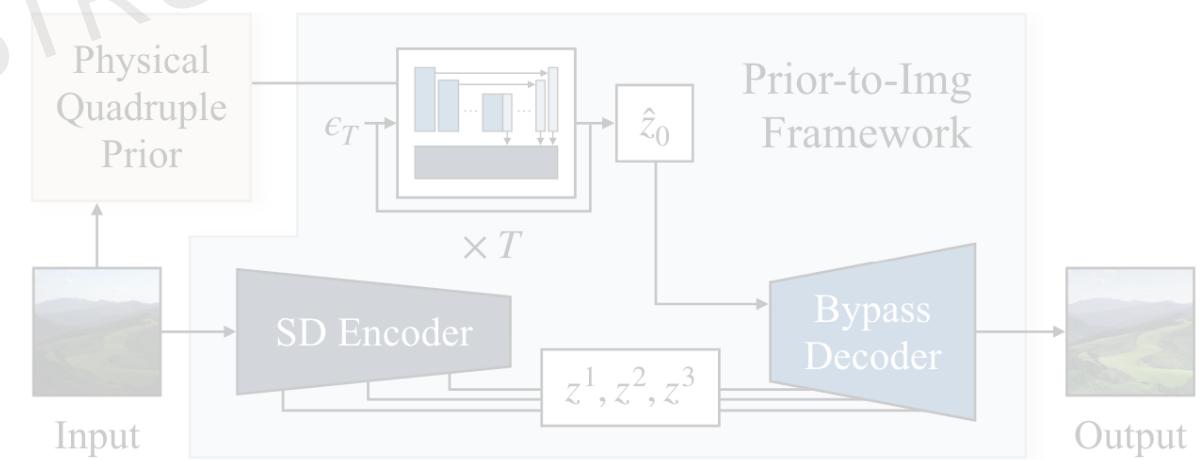
- Improvements to make it more suitable for image restoration task

2. Bypass decoder

- Handle detail distortion through a new decoder

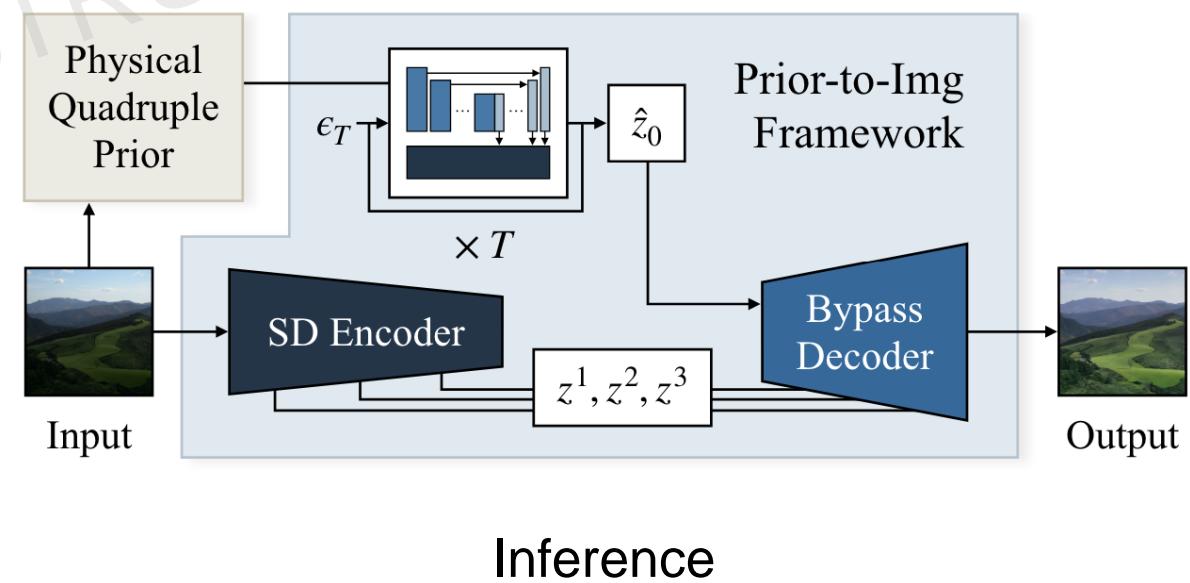
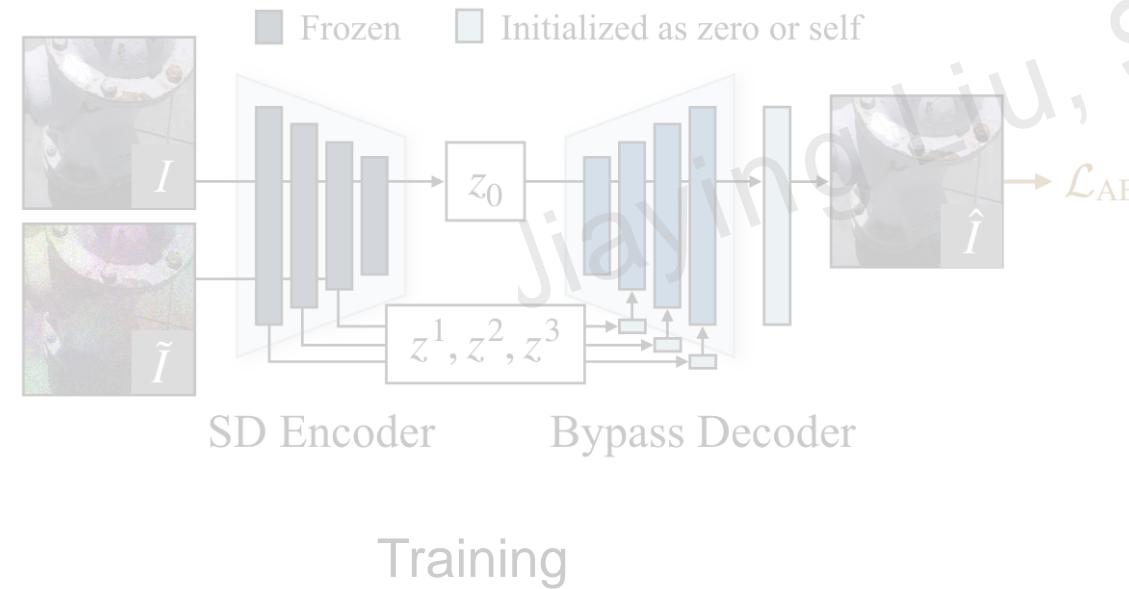


Training



Inference

- Improvements to make it more suitable for image restoration task
 - 2. Bypass decoder
 - Handle detail distortion through a new decoder



Diffusion for Low-Light Enhancement

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- Improvements to make it more suitable for image restoration task
 - Bypass decoder
 - Handle detail distortion through a new decoder



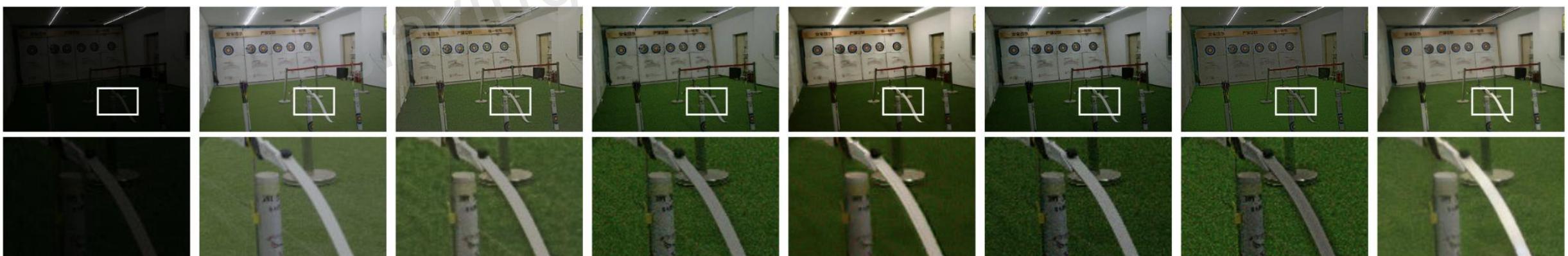
Diffusion for Low-Light Enhancement

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Experiments



Input Ground Truth PairLIE Zero-DCE RUAS-MIT SCI-MIT CLIP-LIT Ours



Input Ground Truth PairLIE Zero-DCE RUAS-LOL SCI-LOL CLIP-LIT Ours

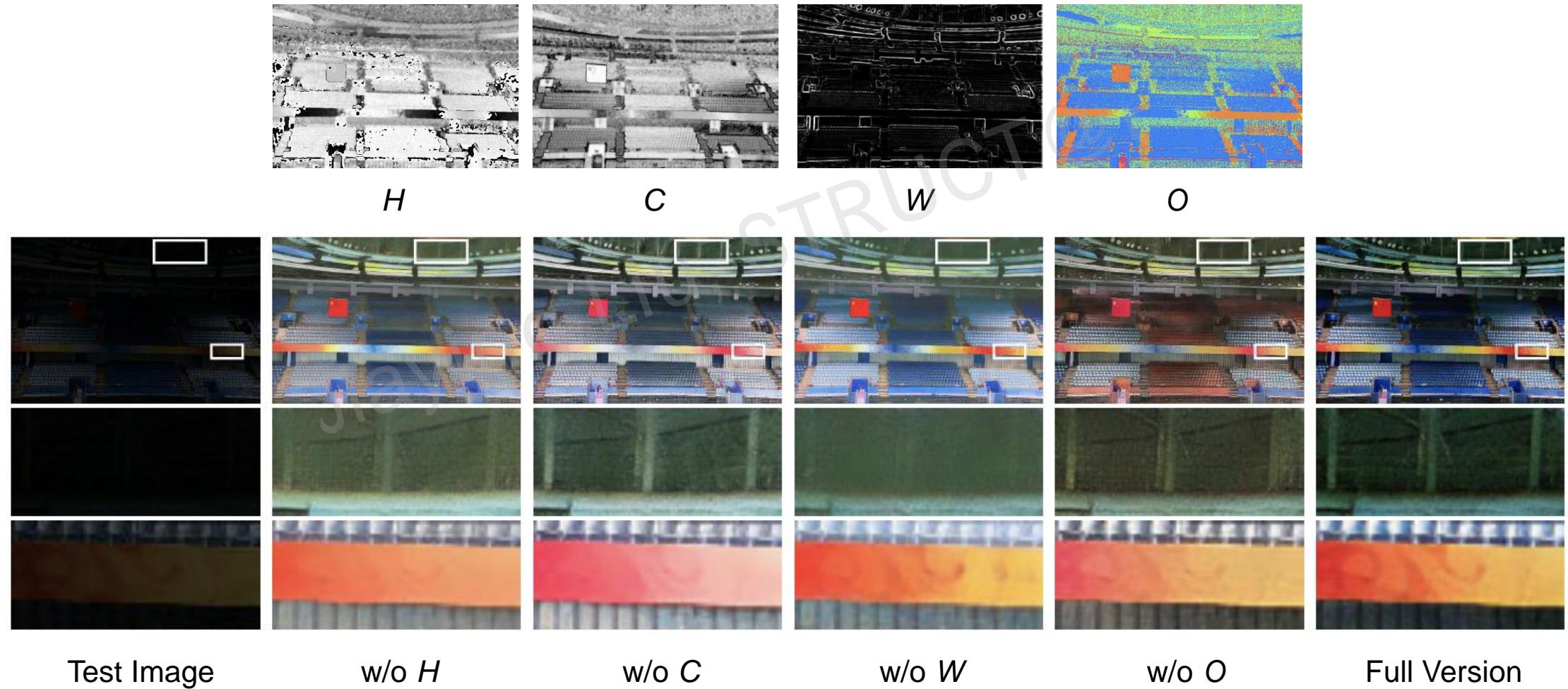
Experiments

| Datasets | | Train Set | LOL [48, 55] | | | | MIT-Adobe FiveK [2] | | | | Unpaired Sets | |
|--------------|-------------------|-------------|--------------|--------------|--------------|--------------|---------------------|--------------|--------|-------|---------------|--------------|
| Metrics | | | PSNR↑ | SSIM↑ | LPIPS↓ | LOE↓ | PSNR↑ | SSIM↑ | LPIPS↓ | LOE↓ | BRISQUE↓ | NL↓ |
| Supervised | Retinex-Net [48] | LOL | 16.19 | 0.403 | 0.534 | 0.346 | 12.30 | 0.687 | 0.258 | 0.244 | 27.10 | 3.254 |
| | KinD [63] | LOL | 20.21 | 0.814 | 0.147 | 0.245 | 14.71 | 0.756 | 0.176 | 0.174 | 26.89 | 0.700 |
| | KinD++ [64] | LOL | 16.64 | 0.662 | 0.410 | 0.288 | 15.76 | 0.650 | 0.319 | 0.176 | 26.16 | 0.431 |
| | URetinex-Net [49] | LOL | 20.93 | 0.854 | 0.104 | 0.245 | 14.10 | 0.734 | 0.182 | 0.187 | 23.80 | 1.319 |
| | Retinexformer [3] | LOL | 28.48 | 0.877 | 0.117 | 0.256 | 13.87 | 0.692 | 0.222 | 0.224 | 14.77 | 1.064 |
| | Retinexformer [3] | MIT | 13.02 | 0.426 | 0.365 | 0.280 | 24.93 | 0.907 | 0.063 | 0.162 | 24.13 | 0.684 |
| | DiffLL [18] | LOL+ | 28.54 | 0.870 | 0.102 | 0.253 | 15.81 | 0.719 | 0.244 | 0.213 | 14.96 | 0.888 |
| Unsupervised | ExCNet [61] | test images | 16.29 | 0.455 | 0.380 | 0.295 | 14.21 | 0.719 | 0.197 | 0.197 | 19.03 | 1.563 |
| | EnlightenGAN [19] | own data | 18.57 | 0.700 | 0.302 | 0.291 | 13.28 | 0.738 | 0.203 | 0.199 | 20.65 | 0.779 |
| | PairLIE [7] | LOL+ | 19.70 | 0.774 | 0.235 | 0.278 | 10.55 | 0.642 | 0.273 | 0.225 | 29.84 | 1.471 |
| | NeRCo [54] | LSRW [12] | 19.67 | 0.720 | 0.266 | 0.310 | 17.33 | 0.767 | 0.208 | 0.213 | 22.81 | 0.603 |
| | CLIP-LIT [27] | own data | 14.82 | 0.524 | 0.371 | 0.320 | 17.00 | 0.781 | 0.159 | 0.194 | 23.44 | 1.962 |
| | ZeroDCE [10] | own data | 17.64 | 0.572 | 0.316 | 0.296 | 13.53 | 0.725 | 0.201 | 0.191 | 21.76 | 1.569 |
| | ZeroDCE++ [23] | own data | 17.03 | 0.445 | 0.314 | 0.391 | 12.33 | 0.408 | 0.280 | 0.417 | 19.34 | 1.150 |
| | RUAS [30] | MIT | 13.62 | 0.462 | 0.346 | 0.292 | 9.53 | 0.610 | 0.301 | 0.272 | 29.91 | 2.091 |
| | RUAS [30] | LOL | 15.47 | 0.490 | 0.305 | 0.330 | 5.15 | 0.373 | 0.669 | 0.399 | 44.70 | 3.312 |
| | RUAS [30] | FACE [56] | 15.05 | 0.456 | 0.371 | 0.292 | 5.00 | 0.366 | 0.685 | 0.398 | 46.21 | 3.633 |
| | SCI [34] | MIT | 11.67 | 0.395 | 0.361 | 0.286 | 16.29 | 0.795 | 0.143 | 0.165 | 16.73 | 0.853 |
| | SCI [34] | LOL+ | 16.97 | 0.532 | 0.312 | 0.289 | 7.83 | 0.573 | 0.360 | 0.187 | 24.46 | 1.893 |
| | SCI [34] | FACE [56] | 16.80 | 0.543 | 0.322 | 0.297 | 10.95 | 0.684 | 0.272 | 0.205 | 18.33 | 1.335 |
| Ours | | COCO [28] | 20.31 | 0.808 | 0.202 | 0.281 | 18.51 | 0.785 | 0.163 | 0.188 | 14.64 | 0.423 |

Diffusion for Low-Light Enhancement

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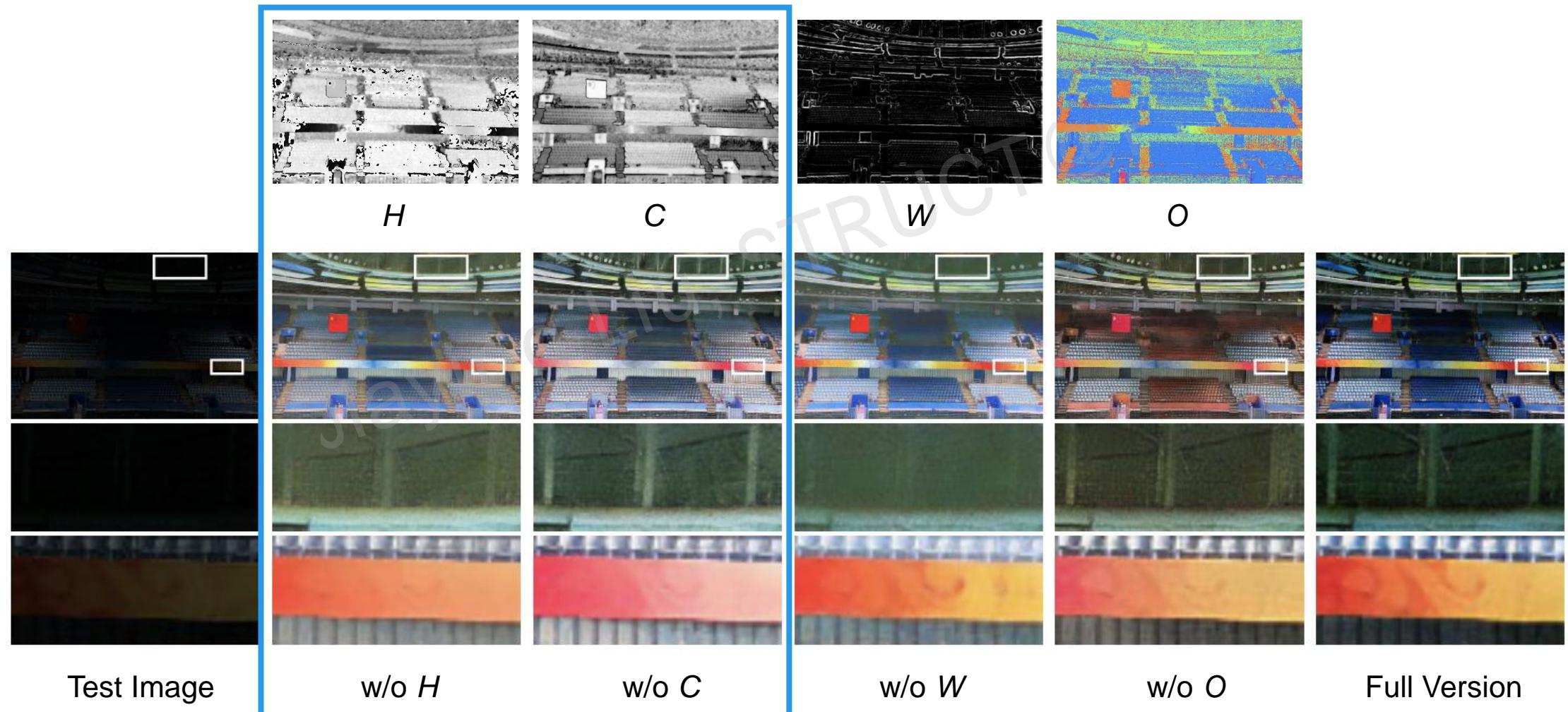
Ablation study on physical quadruple prior



Diffusion for Low-Light Enhancement

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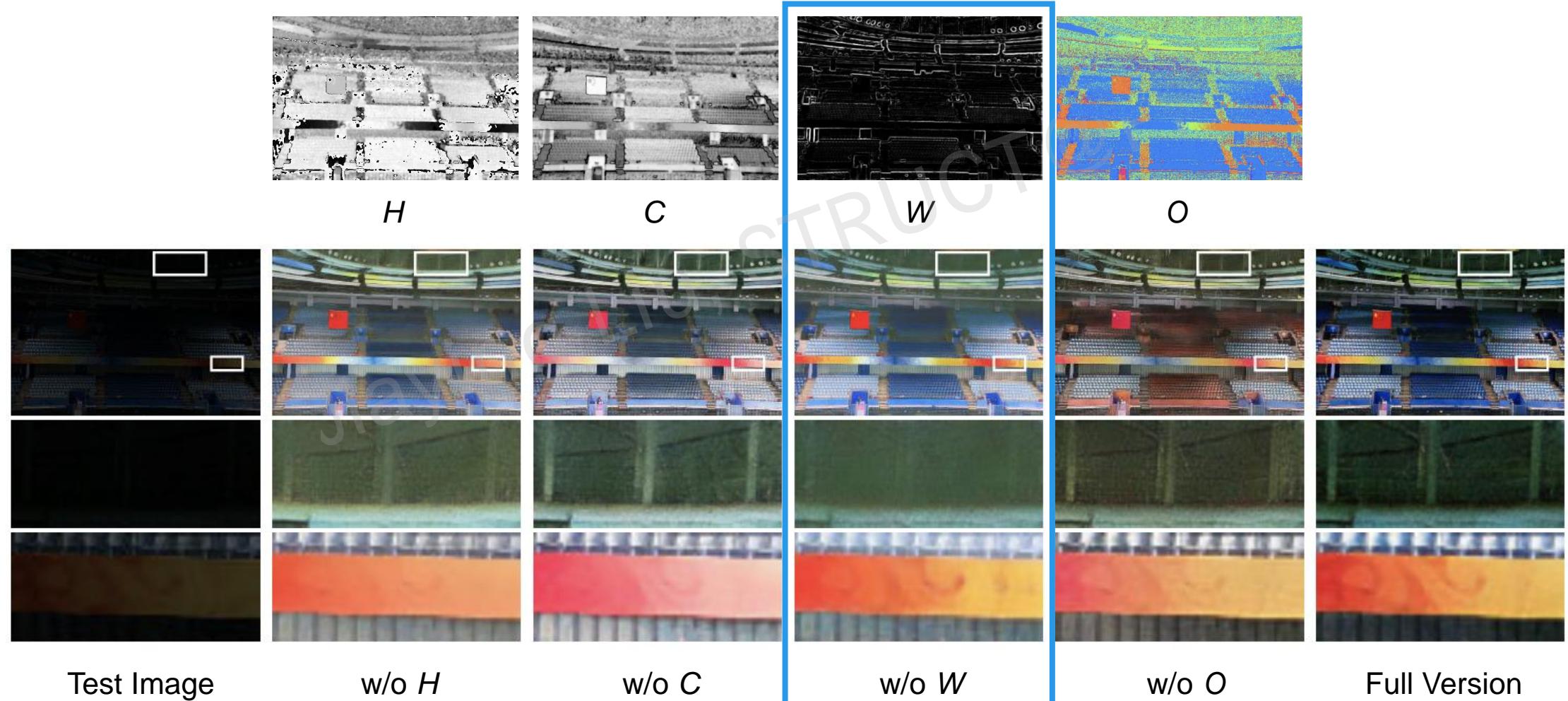
■ Ablation study on physical quadruple prior



Diffusion for Low-Light Enhancement

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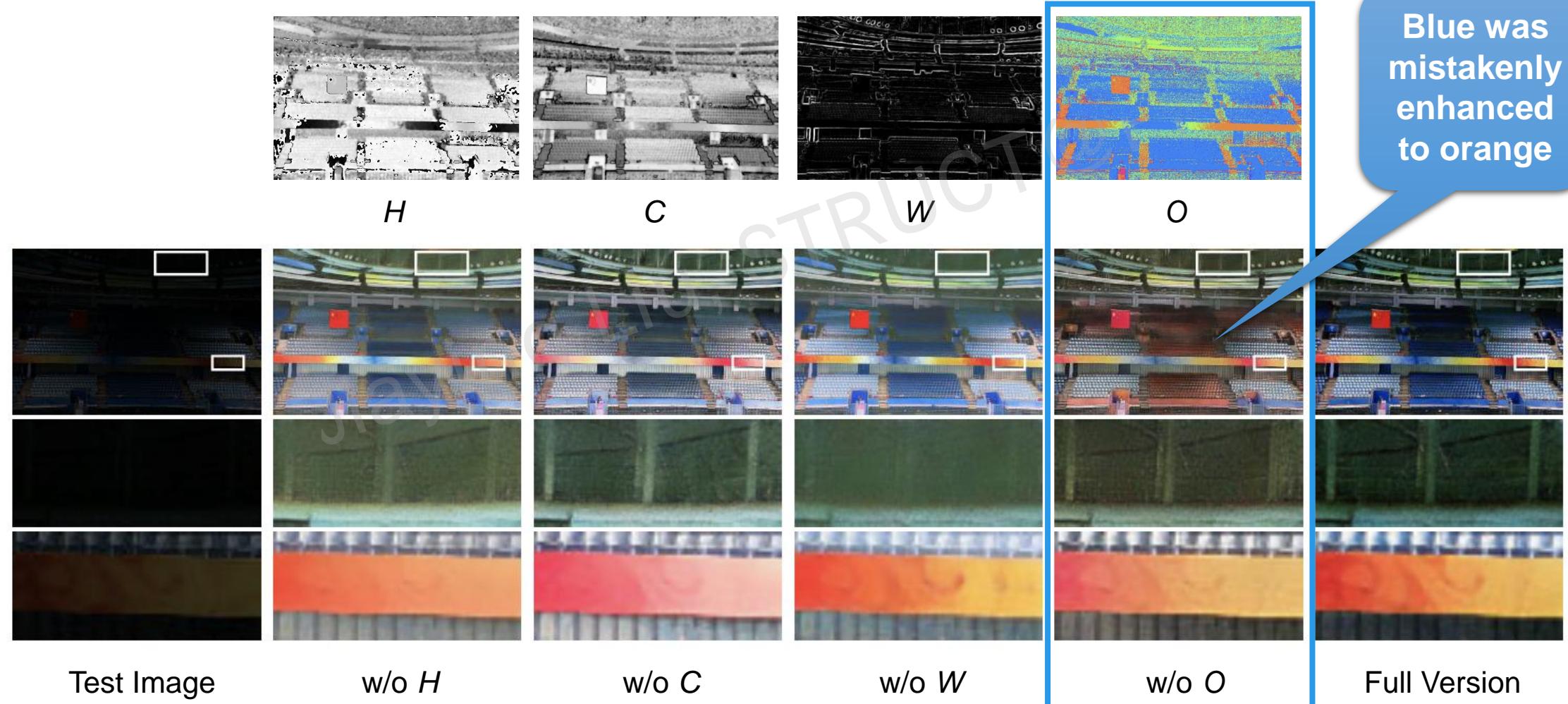
■ Ablation study on physical quadruple prior



Diffusion for Low-Light Enhancement

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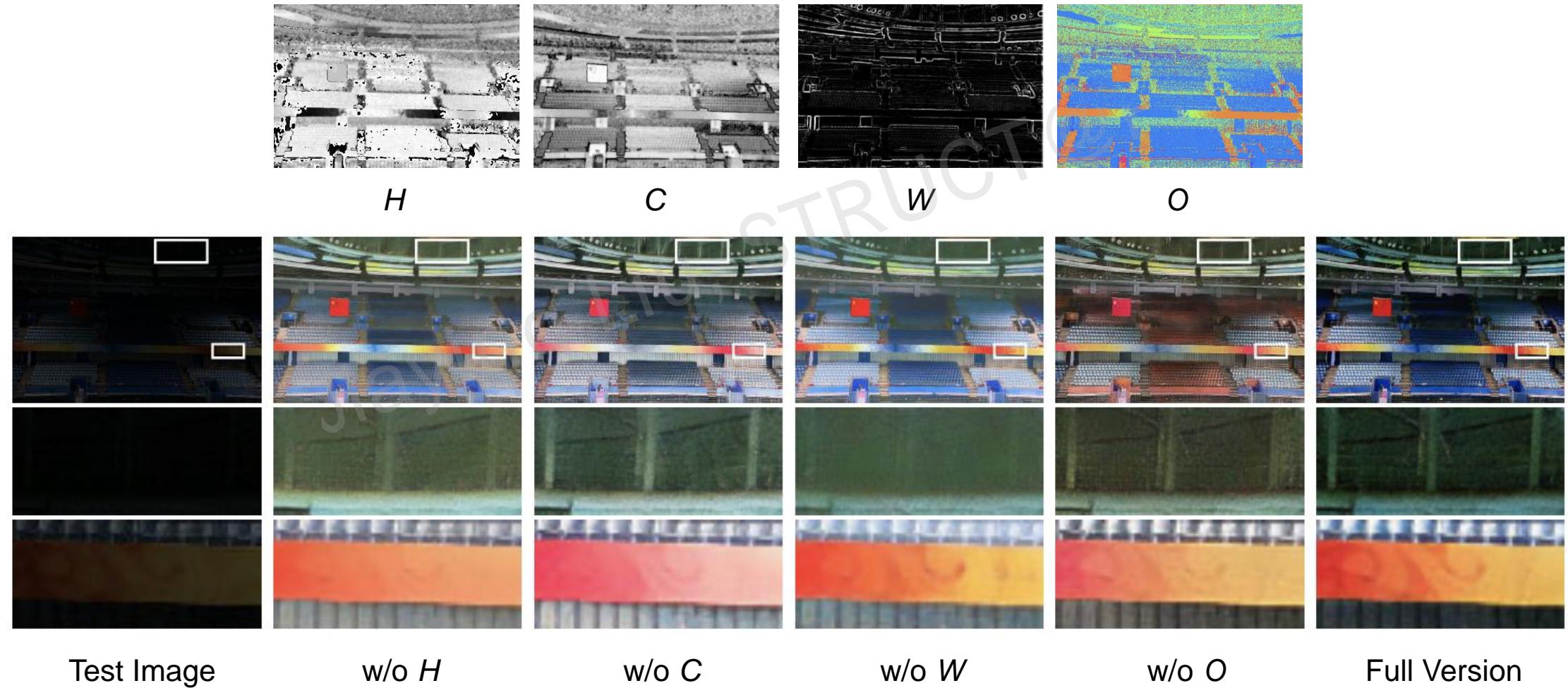
■ Ablation study on physical quadruple prior



Diffusion for Low-Light Enhancement

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■ Ablation study on physical quadruple prior



Diffusion for Low-Light Enhancement

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- Ablation study on prior-to-image mapping



Test Image



SR3 as prior-to-image



Our prior-to-image



Diffusion for Low-Light Enhancement

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Effects of different decoders



Test Image



Stable Diffusion Decoder



Consistency Decoder



Our Bypass Decoder

The proposed decoder can also be applied on the colorization task



Test Image



Stable Diffusion Decoder



Consistency Decoder



Our Bypass Decoder

■ Summary of the method

We propose a **zero-reference** low-light enhancement framework trainable solely with normal light images.

- An **illumination-invariant prior** inspired by physical light transfer.
- A **prior-to-image** framework trained without low-light data
 - Use the prior as a condition to control a pretrained large-scale generative diffusion model
 - Bypass decoder to address the distortion issue



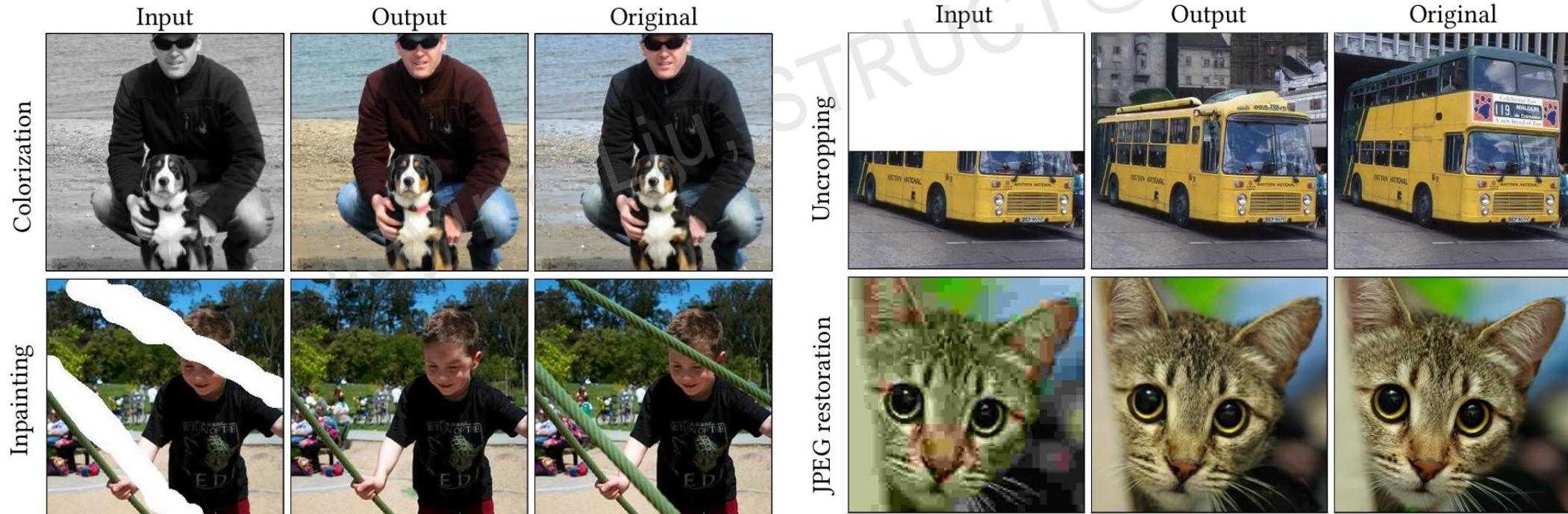
03

Unified Image Restoration and Enhancement

Unified Image Restoration and Enhancement

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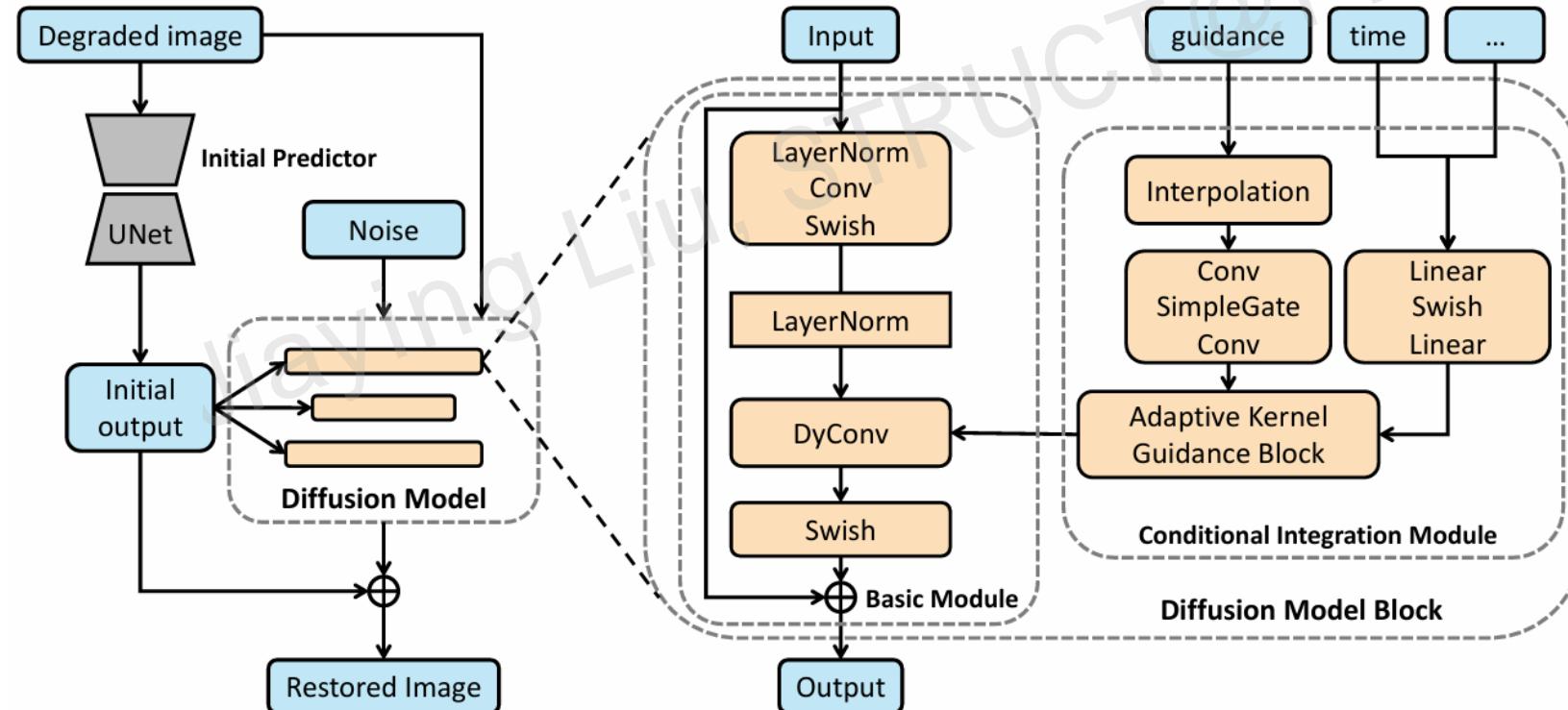
- Palette: Image-to-image diffusion models
 - Based on the SR3 structure: U-Net architecture
 - Different type of conditional images



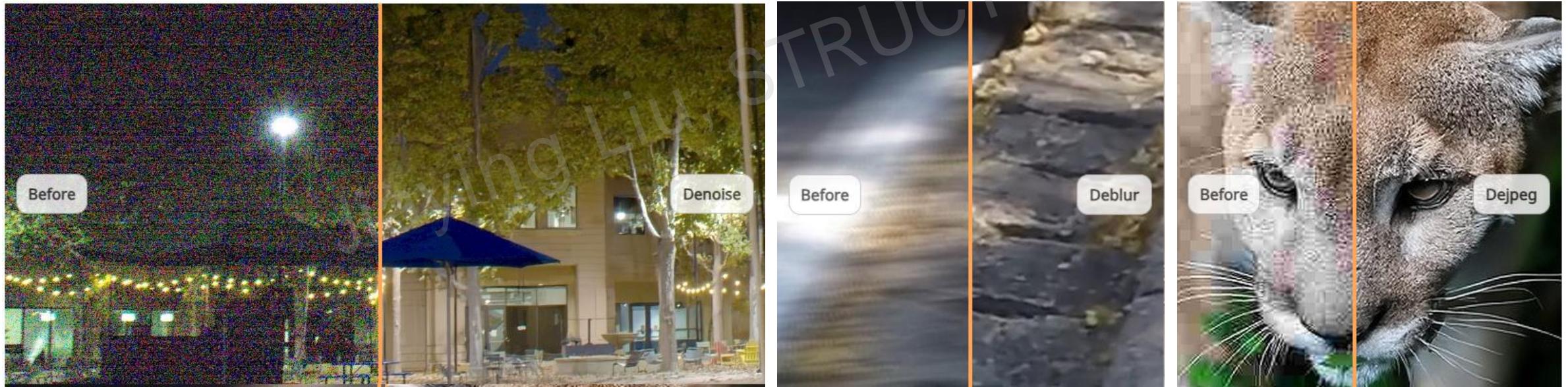
Unified Image Restoration and Enhancement

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- A Unified Conditional Framework for Diffusion-based
 - Same model structure and training hyper-parameters
 - Take different conditional guidance and train on different datasets



- A Unified Conditional Framework for Diffusion-based
 - Same model structure and training hyper-parameters
 - Take different conditional guidance and train on different datasets



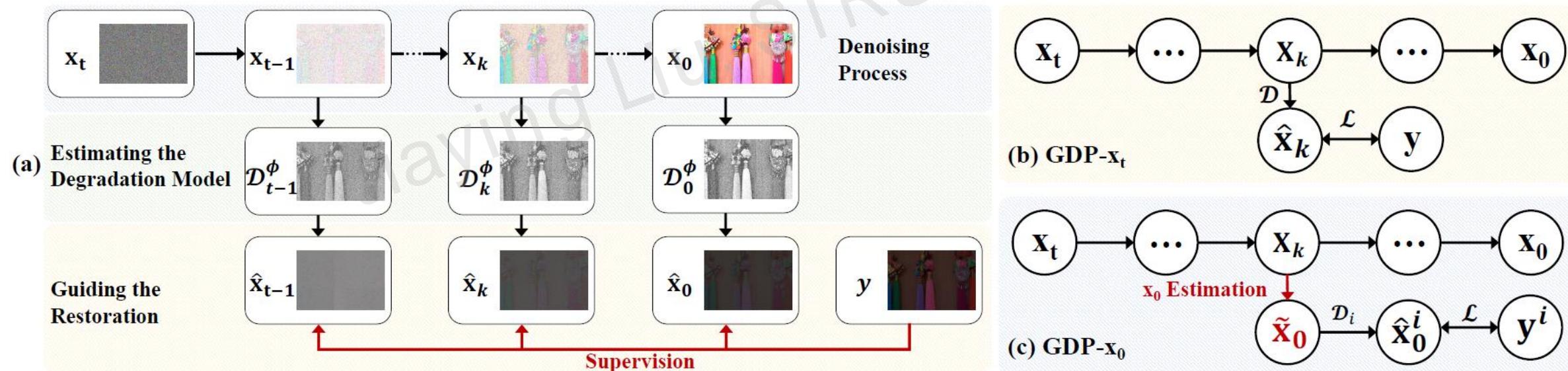
Low-light denoise

Deblur

JPEG restoration

- Generative Diffusion Prior (GDP)

- A unified framework for multiple restoration and enhancement tasks.
- Use a pretrained unconditional image synthesis diffusion model as prior.

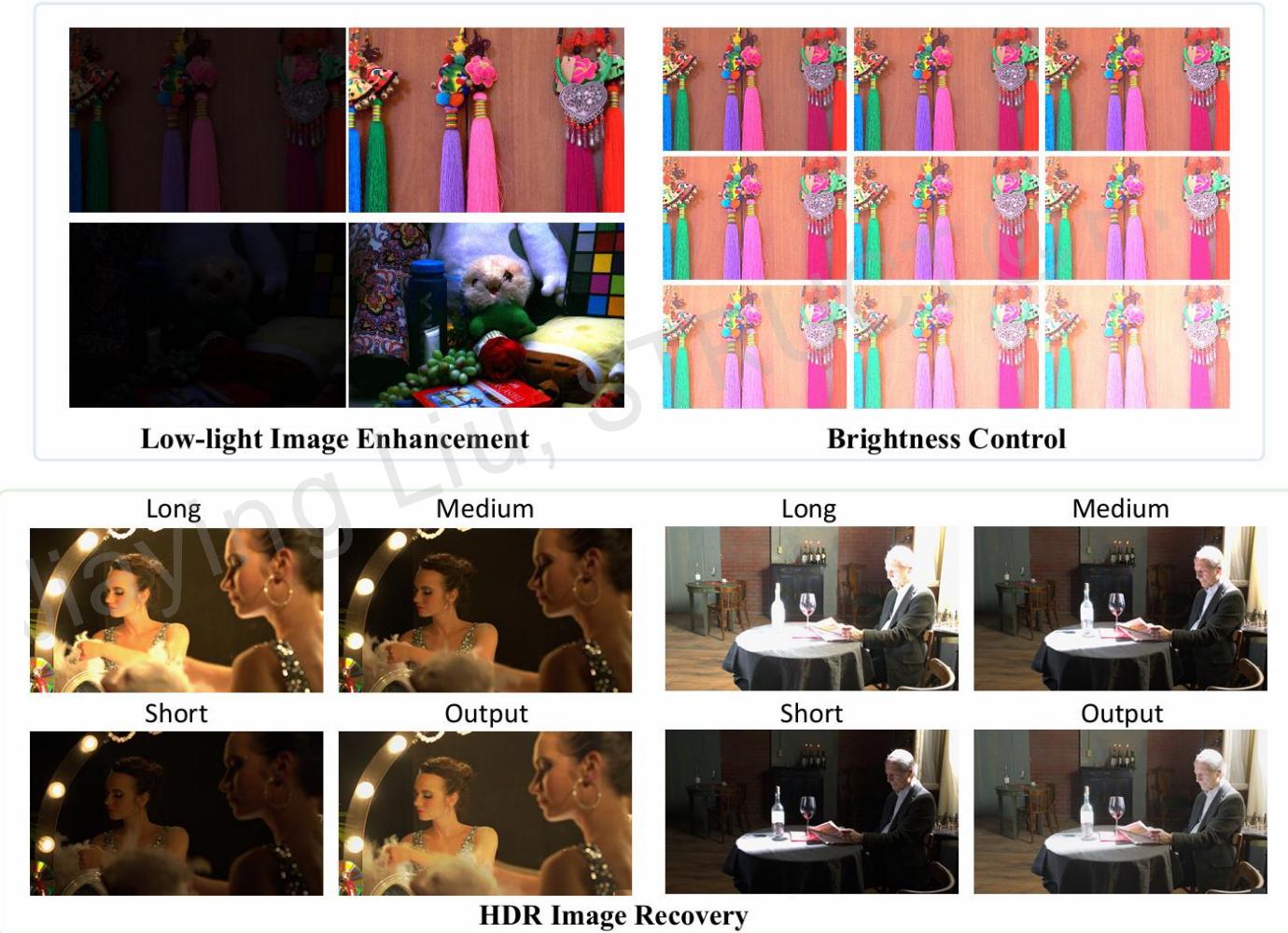


■ Generative Diffusion Prior (GDP)

- A unified framework for multiple restoration and enhancement tasks.
- Use a pretrained unconditional image synthesis diffusion model as prior.
- Different degradation models learned during the sampling process.



■ Generative Diffusion Prior (GDP)





Thank You!

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