The New Era of Image Denoising

Michael Elad

April 10-21, 2023



Computer Science Department The Technion - Israel Institute of Technology Haifa 32000, Israel



Verily Research

A Multiscale Tour of Harmonic Analysis & Machine Learning Celebrating Stéphane Mallat's 60 Birthday



Few Preliminary Words ...

Thank You to Gabriel Peyre and Joan Bruna for inviting me to This Amazing Event!!!



Being here is both a Pleasure and an Honor



Few Preliminary Words ...

Happy Birthday, Stephane and Thank you for your guidance over the years



Invariant Scattering Convolution Networks

Joan Bruna and Stéphane Mallat CMAP, Ecole Polytechnique, Palaiseau, France

Abstract—A wavelet scattering network computes a translation invariant image representation, which is stable to deformations and preserves high frequency information for classification. It cascades wavelet transform convolutions with non-linear modulus and averaging operators. The first network layer outputs SIFT-type descriptors whereas the next layers provide complementary invariant information which improves classification. The mathematical analysis of wavelet scattering networks explain important properties of deep convolution networks for classification. as wavelets. However, wavelet transforms are not invariant to translations. Building invariant representations from wavelet coefficients requires introducing non-linear operators, which leads to a convolution network architecture.

Deep convolution networks have the ability to build large-scale invariants which are stable to deformations er

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This Lecture is About ...



Removal of noise from images is a heavily studied problem in image processing

In this talk we expand on recent discoveries and developments around this seemingly dead topic



Our Agenda

1. Brief Introduction & History

2. Our Focus Today: Denoising for ...

- Image Synthesis
- High perceptual quality recovery

3. Summary



Introduction & History



So, Let's Talk About ...

Image Denoising

or more accurately

Removal of White Additive Gaussian Noise from an Image





Why Assume Gaussian Noise?

□ The Gaussian case is more common and much more important

When considering a Poisson noise,

- High count of photons The distribution gets closer and closer to the Gaussian case
- Low-count Poisson-distributed image can be converted to a Gaussian-noisy one by Anscomb - Variance Stabilizing Transform
- Many of the developed ideas for the Gaussian case can be converted to other noise models
- MMSE denoisers for the
 Gaussian case are of extreme
 theoretical value (see later)







Image Denoising: Evolution





Image Denoising: A Paradigm Shift

How can we design a denoiser?





Image Denoising: Recent Evolution





Our Focus Today: Recent Discoveries



Our Focus Today

Recent findings on using denoisers for other tasks:

Discovery 0: Solving general inverse problems [2013-]

Discovery 1: Image Synthesis [2019-]

Discovery 2: High perceptual quality recovery [2021-]

We turn to describe these results





- In recent years, and with the deep-learning revolution, there is a growing interesting is synthesizing images "out of thin air"
- The popular tool of interest is called GAN Generative Adversarial Network, built of two competing networks – a generator and a critique
- □ Why synthesize? Because
 - We can, and it is fascinating
 - It testifies that we have seized the distribution of images, and
 - It could be used for other needs
- Could we synthesize images differently?



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Question: Given a denoiser $D(y, \sigma)$ how can one synthesize images with it?

Generative modeling by estimating gradients of the data distribution Y Song, S Ermon Advances in Neural Information Processing Systems 32	235	2019
Improved techniques for training score-based generative models Y Song, S Ermon Advances in neural information processing systems 33, 12438-12448	88	2020
Stochastic Solutions for Linear Inverse Problems using the Prior Implicit in a Denoiser Z Kadkhodaie, EP Simoncelli Advances in Neural Information Processing Systems 34	2	2021
nswer: Use $D(y, \sigma)$ as a <i>Projector</i> onto the image	ma	nifol

Practical Implication: Iterated use of $D(\cdot, \sigma)$ with varying σ





Here is the core idea in a nutshell:

Our goal: draw a sample from the distribution of images P(x)

- Start with a random noise image \hat{x}_0
- Climb to a more probable image by the iterative equation:

 $\hat{x}_{k+1} = \hat{x}_k + a \cdot \frac{\nabla \log P(\hat{x}_k)}{P(x_k)} + b \cdot z_k$ (Langevin Dynamics)

This is known (Miyasawa `61) as the Score Function and it is approximately proportional to $[\hat{x}_k - D(\hat{x}_k, \sigma)]$ for a small value of σ

This suggests an implicit relation between MMSE denoisers and Priors: $D(x, \sigma) \leftrightarrow P(x)$

... and this way we got an iterated algorithm that keeps calling to a denoiser, and is guaranteed to obtain a sample from P(x)



In practice, instead of the plain Langevin with a fixed (and small) value of σ we use the Annealed Langevin Algorithm that considers a sequence of blurred priors: Slurred Image

$$\begin{split} P(\mathbf{x} + \mathbf{v}) & \text{for } \mathbf{v} \sim \mathbb{N} \Big(\mathbf{0}, \sigma_k^2 \mathbf{I} \Big) \\ &= P(\mathbf{x}) \otimes \mathbf{c} \cdot \exp \left\{ -\frac{1}{2\sigma^2} \|\mathbf{x}\|^2 \right\} \\ & \text{with } \sigma_0 > \sigma_1 > \sigma_2 \quad \cdots > \sigma_N > 0 \end{split}$$



The core idea: start by drawing from a wider distribution and gradually narrow it, leading to a faster sampling performance



Does it work? Here are some results





Claim: diffusion-based methods are the best in image synthesis









BigGAN (FID 6.95)



Diffusion (FID 4.59)

Diffusion Models Beat GANs on Image Synthesis

Prafulla Dhariwal* OpenAI prafulla@openai.com Alex Nichol* OpenAI alex@openai.com

Abstract

We show that diffusion models can achieve image sample quality superior to the current state-of-the-art generative models. We achieve this on unconditional image synthesis by finding a better architecture through a series of ablations. For conditional image synthesis, we further improve sample quality with classifier guidance: a simple, compute-efficient method for trading off diversity for fidelity using gradients from a classifier. We achieve an FID of 2.97 on ImageNet 128×128 , 4.59 on ImageNet 256×256 , and 7.72 on ImageNet 512×512 , and we match BigGAN-deep even with as few as 25 forward passes per sample, all while maintaining better coverage of the distribution. Finally, we find that classifier guidance combines well with upsampling diffusion models, further improving FID to 3.94 on ImageNet 256×256 and 3.85 on ImageNet 512×512 . We release our code at https://github.com/openai/guided-diffusion.

1 Introduction



Figure 1: Selected samples from our best ImageNet 512×512 model (FID 3.85)



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Surely, you have heard of ...





Surely, you have heard of ...





Suppose that we need to denoise the following image:





Question: How can we denoise an image while targeting "High Perceptual Quality"?

High perceptual quality image denoising with a posterior sampling cgan	3	2021
G Ohayon, T Adrai, G Vaksman, M Elad, P Milanfar		
Proceedings of the IEEE/CVF International Conference on Computer Vision		

Stochastic image denoising by sampling from the posterior distribution52021B Kawar, G Vaksman, M EladProceedings of the IEEE/CVF International Conference on Computer Vision ...

Answer: Denoise by sampling from the posterior P(x|y)

Practical Implication: We consider 2 methods

- Training a deep denoiser via CGAN, or
- Iterated use of an MMSE denoiser $D(\cdot, \sigma)$

These methods produce a multitude of possible solutions



Let's have a closer look at the Stochastic Image Denoiser:

Task: Draw a sample from P(x|y) where $[y = x + n, n \sim \mathbb{N}(0, \sigma_0^2 \mathbf{I})]$

- Start with a random noise image \hat{x}_0
- Climb to a more probable image by the iterative equation:

$$\hat{x}_{k+1} = \hat{x}_k + a \cdot \nabla \log P(\hat{x}_k | y) + b \cdot z_k \qquad \text{Langevin with a conditional Score}$$

$$= \nabla \log P(\hat{x}_k) + \nabla \log P(y | \hat{x}_k)$$
Bayes rule
$$= \hat{x}_k - D(\hat{x}_k, \sigma) + \nabla \log P(y | \hat{x}_k)$$
Approx Score A Gaussian Distribution ?



Let's have a closer look at the Stochastic Image Denoiser:

 $\nabla \log P(\hat{\mathbf{x}}_{k}|\mathbf{y}) = \hat{\mathbf{x}}_{k} - D(\hat{\mathbf{x}}_{k},\sigma) + \nabla \log P(\mathbf{y}|\hat{\mathbf{x}}_{k})$

- As we use the Annealed Langevin algorithm, there are two noise signals to consider:
 - Measurement's noise: $n \sim \mathbb{N}(0, \sigma_0^2 \mathbf{I})$

 $\circ \ \ \text{Synthetic annealing noise:} \ \ v \sim \mathbb{N}\big(0, \sigma_k^2 I\big) \ \text{for} \ \sigma_0 > \sigma_1 > \sigma_2 \ \ \cdots \ > \sigma_N > 0$

Implication: We recover
a sequence of gradually less
noisy images \$\hat{x}_k\$ where their $x_k - D(\hat{x}_k, \sigma_k) + \frac{y - \hat{x}_k}{\sigma_0^2 - \sigma_k^2}$ noise v is assumed to be a portion of n



Stochastic Image Denoiser:

- We start from a noisy image ($\sigma \approx 150$ in this example)
- Then gradually denoise it using (conditional) annealed Langevin dynamics







Let's have a closer look at the Conditional GAN Denoiser:



- Typical design approach: Optimize a distortion measure (e.g. MSE) between the denoised and the ideal images
- Adversarial loss could be added to improve the perceptual quality





However, when used together, we get a compromise between distortion and perceptual quality





- □ For ill-posed restoration tasks, perceptual quality performance comes at the expense of its distortion [Blau & Michaeli 2017]
- We aim for best perceptual quality
- The posterior distribution attains perfect perceptual quality, compromising 3dB compared to the MMSE [Blau & Michaeli 2017]
- We propose to sample from the posterior via a Conditional GAN mechanism (PSCGAN)



Samples from $P_{X|Y=y}$

$$x \sim P_X$$
 $y \sim P_{Y|X=x}$





The PSCGAN Architecture:





What about the Loss?

CGAN optimization leads to posterior sampling [Adler et al. 2018]:

 $\overline{\min_{\theta} \max_{\omega} \mathbb{E}_{X,Y}} \left[C_{\omega}(x,y) \right] - \mathbb{E}_{D_{\theta},Y,Z} \left[C_{\omega}(D_{\theta},y) \right]$

However, this requires an unavailable balanced dataset to train on (many x's for each y and many y's for each x)

On the other hand, we would like to avoid a penalty of the form

 $\mathbb{E}_{X,Y,Z} \big[\| x - D_{\theta}(y,z) \|_2^2 \big]$

Our remedy: adding an MMSE oriented penalty term:

$\mathbb{E}_{X,Y} \big[\| \mathbf{x} - \mathbb{E}_{\mathbf{z}} [\mathbf{D}_{\theta} | \mathbf{y}] \|_2^2 \big]$

This gives the MMSE result "for free" (averaging many instances)



CGAN:



Oh ... and One Last Thing



Goal: Recovery from corrupted measurements

De-Noising De-Blurring In-Painting **De-Mosaicing** Tomography Image Scale-Up & super-resolution

Can we suggest a "sampler" from P(x|y) for handling all these problems, in order to obtain "perfect looking" results?

Answer: Yes! Use Langevin dynamics again, in an adapted form







= Hx + n



□ However, ... while $n - Hv_i$ is a simple Gaussian, it's dependency on x_i in non-trivial, so how do we proceed from here?



 \Box Step 1: Use SVD for decoupling the measurements $H = U\Sigma V^{T}$:

$$U^{T}y = U^{T}[U\Sigma V^{T}(x_{i} - v_{i}) + n] = \Sigma V^{T}(x_{i} - v_{i}) + U^{T}n$$

$$\underbrace{ \bigvee y = Hx + n}_{}$$

$$y_{T}[k] = \sigma_{k}\tilde{x}_{T}[k] - \sigma_{k}\tilde{v}_{T}[k] + n_{T}[k]$$

Decouple $\tilde{x}_T[k] \leftrightarrow \tilde{v}_T[k]$ by choosing $\tilde{v}_T[k]$ to be a portion of $n_T[k]$

 $\hfill Thus, we can apply the Langevin dynamics algorithm on <math display="block">\widetilde{x}_T = V^T x_i \text{ given } y_T = U^T y \text{ and sample from the conditional}$

Bottom line: An MMSE denoiser is used for a novel solution of inverse problems, this time targeting best perceptual quality



Noisy Inpainting: A portion missing and noise with $\sigma_0 \approx 25$





Super resolution: downscaling by 4 with additive noise of $\sigma_0 \approx 25$





Super resolution: downscaling by 4 with additive noise of $\sigma_0 \approx 12$



Original

Degraded



Deblurring: uniform 5 \times 5 blur with additive noise of $\sigma_0 \approx 25$





Compressive sensing (12.5%) with additive noise of $\sigma_0 \approx 25$





And just to remind you ...

The proposed diffusion-based sampling scheme, while quite appealing, suffers from several key shortcomings:

It is rather SLOW (many denoising activations)
 It is limited (as of now) to specific families of images
 Relying on SVD is cumbersome

See our recent work that answers (some of) these challenges:

Denoising Diffusion Restoration Models B Kawar, M Elad, S Ermon, J Song Advances in Neural Information Processing Systems (NeurIPS)



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2022

Time to Summarize



What Have we Seen Today?

Suppose that we are given an MMSE denoiser D(y)





Summary

Image Denoising ... Not What You Think

- There are so many opportunities and challenges in better understanding, designing, and proposing creative usage of image denoisers
- Despite the fact that this has not been a talk about Deep-Learning, the presence of this field in the topics covered is prominent

