DEEP PROBABILISTIC SUBSAMPLING FOR TASK-ADAPTIVE COMPRESSED SENSING

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ABSTRACT

The field of deep learning is commonly concerned with optimizing predictive models using large pre-acquired datasets of densely sampled datapoints or signals. In this work, we demonstrate that the deep learning paradigm can be extended to incorporate a subsampling scheme that is jointly optimized under a desired minimum sample rate. We present Deep Probabilistic Subsampling (DPS), a widely applicable framework for task-adaptive compressed sensing that enables end-to-end optimization of an optimal subset of signal samples with a subsequent model that performs a required task. We demonstrate strong performance on reconstruction and classification tasks of a toy dataset, MNIST, and CIFAR10 under stringent subsampling rates in both the pixel and the spatial frequency domain. Due to the task-agnostic nature of the framework, DPS is directly applicable to all real-world domains that benefit from sample rate reduction.

1 Introduction

In many real-world prediction problems, acquiring data is expensive and often bandwidth constrained. Such is the case in regimes as medical imaging (Lustig et al., 2007; Choi et al., 2010; Chernyakova & Eldar, 2014), radar (Baraniuk, 2007), and seismic surveying Herrmann et al. (2012). By carefully reducing the number of samples acquired over time, in pixel-coordinate-space or in k-space, efficient subsampling schemes lead to meaningful reductions in acquisition time, radiation exposure, battery drain, and data transfer.

Subsampling is traditionally approached by exploiting expert knowledge of the signal of interest. Famously, the Nyquist theorem states that when the maximum frequency of a continuous signal is known, perfect reconstruction is possible when sampled at twice this frequency. More recently, it has been shown that if the signal is known to be sparse in some domain, sub-Nyquist rate sampling can be achieved through compressive measurements and subsequent optimization of a linear system under said sparsity prior; a framework known as compressed sensing (CS) (Donoho et al., 2006; Eldar & Kutyniok, 2012; Baraniuk, 2007).

These methods, however, are lacking in the sense that they do not fully exploit both the underlying data distribution and the information required to solve the downstream task of interest, such as disease prediction or semantic segmentation. Formalizing such knowledge is challenging in its own right and would require careful analysis for each modality and downstream task. In this work, we propose to explore the deep learning hypothesis as a promising alternative: reducing the need for expert knowledge in lieu of large data-sets and end-to-end optimization of neural networks.

As subsampling is non-differentiable, incorporating it into an end-to-end optimized deep learning model is non-trivial. Here, we take a probabilistic approach: rather than learning a subsampling scheme directly, we pose a probability distribution that expresses belief over effective subsampling patterns and optimize the distribution's parameters instead. To enable differentiable sampling from this distribution, we leverage recent advancements in continuous relaxation of the categorical distribution, known as the Gumbel-softmax or Concrete distribution (Jang et al., 2017; Maddison et al., 2016). This enables end-to-end training of both the subsampling scheme and the downstream model.

Naively, the number of parameters of a distribution over an n-choose-k problem scales factorially, which is intractable for all practical purposes. When it comes to subset sampling, recent work (Kool et al., 2019; Xie & Ermon, 2019) has explored a more tractable yet limited fully-factorized

parametrization with a single parameter per sample, for which the top-K samples are stochastically selected. Although efficient, for the task of subsampling this aggressive factorization can not express the scenario where multiple samples can be equally good, yet redundant in combination, and forces a model to focus on a single mode right off the bat.

We propose a novel parameterization for the subsampling distribution by conditioning on the output sample index. This leads to an expressive yet tractable distribution that prevents redundant sampling, whilst allowing equal weight attribution to samples that are equally informative in isolation. This is essential for the efficient exploration of potential subsampling schemes, as it enables the model to maintain multiple hypotheses during optimization.

Our main contributions are as follows:

- DPS: A new regime for task-adaptive subsampling using a novel probabilistic deep learning framework for jointly learning a sub-Nyquist sampling scheme with a predictive model for downstream tasks.
- A novel parametrization of the subsampling distribution by conditioning on the output sample index, balancing tractability and exploration and outperforming fully factorized parametrizations.
- Improved performance over strong subsampling baselines in image classification and reconstruction in both Fourier and pixel space.

2 Related work

Some recent works have proposed deep-learning-based subsampling methods for fast MRI, historically being one of the most prominent applications of CS. Weiss et al. (2019) exploit gradient backpropagation to a fixed set of real-valued coordinates, enabled by their subsequent (limited-support) interpolation on the discrete k-space grid, and Bahadir et al. (2019) formulate the sampling problem by learning pixel-based thresholding of i.i.d. samples drawn from a uniform distribution. Where the former suffers from limited exploratory capabilities (likely due to its compact support), with learned sampling schemes typically not deviating far from their initialization, the latter controls the sample rate only indirectly, through the addition of a sparsity-promoting ℓ_1 penalty on the mask.

Closely related to our methodology, Xie & Ermon (2019); Kool et al. (2019); Plötz & Roth (2018) leverage an extension to the Gumbel-max trick (Gumbel, 1954) for subset selection using a categorical distribution with N-1 free parameters. They rely on relaxed top-K sampling, as first proposed by Vieira (2014). These methods build upon the continuous relaxation of the categorical distribution known in the deep learning community as the Gumbel-softmax trick or the concrete distribution (Jang et al., 2017; Maddison et al., 2016). Our method builds upon the same foundation, but proposes to parameterize the sub-sampling scheme conditioned on the output sample index, leading to $N-1\times K$ free parameters for N input samples and K selected samples.

We differentiate our contribution from deep encoder-decoder methods for data compression (Baldi & Hornik, 1989; Hinton & Zemel, 1993; Blier & Ollivier, 2018; Habibian et al., 2019), which do not aim at reducing data rates already at the sensing and digitization stage. Related work by Mousavi et al. (2019) and Wu et al. (2019), focusses on the problem of learning compressive linear encoders/filters, rather than discrete subsampling as addressed here.

Through the lens of contemporary deep learning, subsampling can be interpreted as a form of *attention* (Bahdanau et al., 2014; Kim et al., 2017; Parikh et al., 2016; Vaswani et al., 2017). Rather than attending on intermediate representations, our model "attends" directly on the input signal. For sub-sampling to be effective, sparse weights are essential. In the space of attention, this is known as *hard* attention (Xu et al., 2015), and is typically optimized using the REINFORCE gradient estimator (Williams, 1992). In contrast to the method of attention as applied in these works, our method aims for a a fixed, reduced subsampling rate.

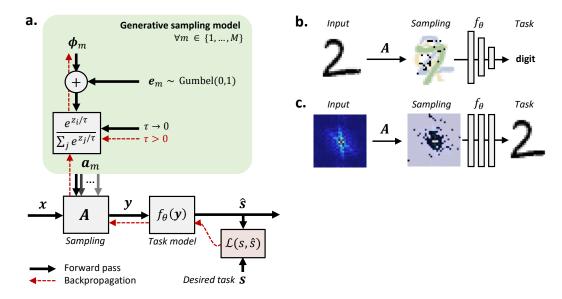


Figure 1: (a) System-level overview of the proposed framework, in which a probabilistic generative sampling model and a subsequent task model are jointly trained to fulfil a desired system task. (b,c) Two illustrative task-based sampling paradigms: image classification from a partial set of pixels (b), and image reconstruction from partial Fourier measurements (c), respectively.

3 Method

3.1 TASK-ADAPTIVE SYSTEM MODEL

We consider the problem of performing some downstream task s through a learned subsampling scheme A, resulting in measurements $\mathbf{y} \in \mathbb{R}^M$ of an underlying fully-sampled signal $\mathbf{x} \in \mathbb{R}^N$, with M << N:

$$\mathbf{y} = \mathbf{A}\mathbf{x},\tag{1}$$

where $\mathbf{A} \in \{0,1\}^{M \times N}$ is the subsampling measurement matrix. We here concern ourselves specifically with scenarios in which the rows of \mathbf{A} are constrained to having cardinality one, i.e. $||a_m||_0 = 1, \forall m \in \{0,...,M-1\}$. In such cases, \mathbf{A} serves as a subset selector, sampling M out of N elements in \mathbf{x} . From the resulting low-rate measurements \mathbf{y} we aim at performing task s through:

$$\hat{s} = f_{\theta}(\mathbf{y}),\tag{2}$$

with $f_{\theta}(\mathbf{y})$ being a function that is differentiable with respect to its input and parameters θ , e.g. a neural network.

Given a downstream task and dataset, we are interested in learning both the optimal processing parameters θ , and the sampling scheme as described in eq. (1). To circumvent the non-differential nature of discrete sampling, we will introduce a novel fully probabilistic sampling strategy that allows for gradient-based learning through error backpropagation, on which we detail in the following section.

3.2 Trainable probabilistic sampling

Figure 1 shows a schematic overview of the proposed framework. Since direct optimization of the elements in $\bf A$ is intractable due to its combinatorial nature, we here instead propose to leverage a tractable generative sampling model that is governed by a learned subsampling distribution, parameterized by $\bf \Phi$:

$$\mathbf{A} \sim P(\mathbf{A}|\mathbf{\Phi}). \tag{3}$$

Thus, rather than optimizing A, we optimize the distribution parameters Φ . To warrant sufficient expressiveness while maintaining tractability, we learn the parameters $\phi_m \in \mathbb{R}^N$ of M independent

categorical distributions (rather than the joint distribution, which scales factorialy), being the rows of $\Phi \in \mathbb{R}^{M \times N}$.

Formally, we define each m^{th} measurement row $\mathbf{a}_m \in \{0,1\}^N$ in \mathbf{A} as a one-hot encoding of an independent categorical random variable $\mathbf{r}_m \sim \operatorname{Cat}(N, \boldsymbol{\pi}_m)$. We define $\boldsymbol{\pi}_m \in \mathbb{R}^N = \{\pi_{m,1}, \dots, \pi_{m,N}\}$, being a vector containing N class probabilities, and parameterize it in terms of its unnormalized logits $\phi_{m,n}$ such that:

$$\pi_{m,n} = \frac{\exp \phi_{m,n}}{\sum_{i=1}^{N} \exp \phi_{m,i}}.$$
 (4)

We sample from $\operatorname{Cat}(N, \pi_m)$ by leveraging the Gumbel-max trick (Gumbel, 1954), a reparameterization of the sampling procedure that is a function of the distribution parameters and a Gumbel noise vector $\mathbf{e}_m \in \mathbb{R}^N$ with i.i.d. Gumbel noise samples $e_{m,n} \sim \operatorname{Gumbel}(0,1)$. A realization $\tilde{\mathbf{r}}_m$ is then defined as:

$$\tilde{\mathbf{r}}_m = \underset{n}{\operatorname{argmax}} \{ \operatorname{WR} (\phi_{m,n} + e_{m,n}) \}, \quad m \in \{0, \dots, M-1\}.$$
 (5)

Operator WR indicates sampling Without Replacement, which we implement by evaluating eq. (5) sequentially, masking previous realizations $\tilde{r}_1, \dots, \tilde{r}_{m-1}$. Introducing the function one_hot $_N(\cdot)$ as the operator that returns a one-hot vector of length N, we finally obtain:

$$\mathbf{a}_{m} = \operatorname{one_hot}_{N} \left\{ \tilde{\mathbf{r}}_{m} \right\} =$$

$$= \operatorname{one_hot}_{N} \left\{ \operatorname{argmax}_{n} \left\{ \operatorname{WR} \left(\phi_{m,n} + e_{m,n} \right) \right\} \right\}. \tag{6}$$

To permit error backpropagation for efficient optimization of Φ , we require $\nabla_{\phi_m} \mathbf{a}_m$ to exist $\forall m \in \{1, \dots, M\}$. Since $\operatorname{argmax}(\cdot)$ is a non-differentiable operator, we adopt the Straight-Through Gumbel Estimator (Jang et al., 2017; Maddison et al., 2016) as a surrogate for $\nabla_{\phi_m} \mathbf{a}_m$:

$$\nabla_{\boldsymbol{\phi}_{m}} \mathbf{a}_{m} := \nabla_{\boldsymbol{\phi}_{m}} \mathbb{E}_{\mathbf{e}_{m}} \left[\operatorname{softmax}_{\tau} (\operatorname{WR} (\boldsymbol{\phi}_{m} + \mathbf{e}_{m})) \right] =$$

$$\nabla_{\boldsymbol{\phi}_{m}} \mathbb{E}_{\mathbf{e}_{m}} \left[\frac{\exp\{(\boldsymbol{\phi}_{m} + \mathbf{e}_{m})/\tau\}}{\sum_{i=1}^{N} \exp\{(\boldsymbol{\phi}_{m,i} + e_{m,i})/\tau\}} \right], \tag{7}$$

with (row operator) softmax_{τ}(·) as a continuous differentiable approximation of the one-hot encoded $\operatorname{argmax}(\cdot)$ operation. See appendix A for the full derivation of $\nabla_{\phi_m} \mathbf{a}_m$.

We refer to sampling using the $\operatorname{softmax}_{\tau}(\cdot)$ function as soft sampling. Its temperature parameter τ serves as a gradient distributor over multiple entries (i.e. logits) in ϕ_m . Using a relatively high value enables updating of multiple logits during training, even though a hard sample was taken in the forward pass. In the limit of $\tau \to 0$, soft sampling approaches the one-hot encoded $\operatorname{argmax}(\cdot)$ operator in eq. (6) (Jang et al., 2017; Maddison et al., 2016). As such, we define:

$$\mathbf{a}_m \coloneqq \lim_{\tau \to 0} \operatorname{softmax}_{\tau}(\operatorname{WR}(\phi_m + \mathbf{e}_m)), \text{ and}$$
 (8)

$$\nabla_{\boldsymbol{\phi}_{m}} \mathbf{a}_{m} := \nabla_{\boldsymbol{\phi}_{m}} \mathbb{E}_{\mathbf{e}_{m}} \left[\operatorname{softmax}_{\tau} (\operatorname{WR} \left(\boldsymbol{\phi}_{m} + \mathbf{e}_{m} \right)) \right], \quad \tau > 0,$$
(9)

with $m \in \{1, ..., M\}$.

4 EXPERIMENTS

We test the applicability of the proposed task-adaptive DPS framework for three datasets and two distinct tasks: image classification and image reconstruction. We explore subsampling in pixel coordinate space as well as in k-space. The latter is relevant for scenarios in which data is acquired in the frequency domain, such as (but not limited to) magnetic resonance imaging (MRI).

4.1 MNIST CLASSIFICATION

Experiment setup Classification performance was tested on the MNIST database (LeCun et al., 1998), comprising $70,000\ 28\times28$ grayscale images of handwritten digits 0 to 9. We split the dataset into 50,000 training images, 5,000 validation, and 5,000 test images. We train our sampling

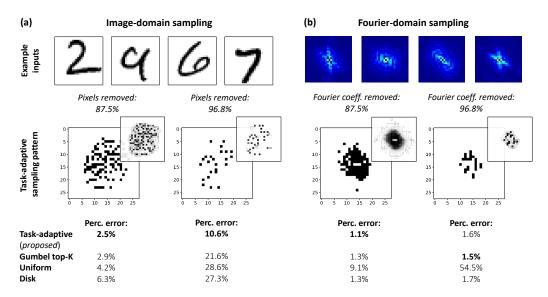


Figure 2: MNIST classification for (a) image-domain, and (b) Fourier-domain subsampling. (top) Several example images in the respective sampling domains, (middle) learned task-adaptive subsampling patterns, with their relative sample incidence across a 1000 such realizations (inset), and (bottom) classification results of the proposed task-adaptive scheme compared to Gumbel top-K subset learning and two non-learned baseline sampling approaches.

model to take partial measurements in either the image or Fourier domain image, and process them through the task model to yield a classification outcome. We compare our results to those obtained using uniformly distributed pixel/Fourier samples, a sampled disk/low-pass filter, and Gumbel top-K sampling (Kool et al., 2019).

Task model After sampling M elements, all N zero-masked samples (or 2N in the case of complex Fourier samples) are passed through a series of 5 fully-connected layers, having N, 256, 128, 128 and 10 output nodes, respectively. The activations for all but the last layer were leaky ReLUs, and 40% dropout was applied after the third layer. The 10 outputs were normalized by a Softmax function to yield the respective classification probabilities. Zero-filling and connecting all possible samples (rather than only connecting the M selected samples) facilitated faster co-adaptation of the network to different sampling patterns during training.

Training details We train the network to maximize the log-likelihood of the observations $\mathcal{D} = \{(\mathbf{x}_i, \mathbf{s}_i) \mid i \in 0, \dots, L\}$ through minimization of the categorical cross-entropy between the predictions and the labels, denoted by \mathcal{L}_s . We moreover promote training towards one-hot sampling distributions $\boldsymbol{\pi}_m$ by penalizing high entropy:

$$\mathcal{L}_e = -\sum_{m=1}^{M} \sum_{n=1}^{N} \pi_{m,n} \log \pi_{m,n}, \tag{10}$$

with $\pi_{m,n}$ defined as in eq. (4). The total optimization problem is thus:

$$\hat{\mathbf{\Phi}}, \hat{\theta} = \underset{\mathbf{\Phi}, \theta}{\operatorname{argmin}} \Big\{ \mathbb{E}_{(\mathbf{x}, \mathbf{s}) \sim p_{\mathcal{D}}} \mathcal{L}_s + \mu \mathcal{L}_e \Big\}, \tag{11}$$

where $p_{\scriptscriptstyle D}$ is the data generating distribution and $\mu=1e-6$. The temperature parameter τ in eq. (7) was set to 10, and the sampling distribution parameters Φ were initialized randomly, following a zero-mean Gaussian distribution with standard deviation 0.25. Equation 11 was optimized using stochastic gradient descent on batches of 32 examples, approximating the expectation by a mean across the train dataset. To that end, we used the ADAM solver ($\beta_1=0.9,\ \beta_2=0.999,\$ and $\epsilon=1e-7$) (Kingma & Ba, 2014), training for 150 epochs. We adopted different learning rates for the sampling parameters Φ and the parameters of the task model θ , being 1e-3 and 1e-4, respectively.

Results The results presented in Figure 2a show that task-adaptive image-domain sampling significantly outperforms the fixed sampling baselines, as well as the alternative Gumbel top-K-based subset selector (Kool et al., 2019). The resulting patterns qualitatively demonstrate how, for this task, a sensible selection of pixels that are most informative was made (slightly slanted, and capturing discriminative areas). Notably, partial Fourier measurements (Figure 2b) allowed for a much greater reduction of the number samples, with task-adaptive sampling again outperforming the fixed sampling approaches. Interestingly, the DC and very-low frequency components were consistently not selected.

4.2 'LINES AND CIRCLES' IMAGE RECONSTRUCTION

Experiment setup To evaluate reconstruction of structured images from highly undersampled partial Fourier (k-space) measurements (keeping 3.1% of the coefficients), we generated synthetic toy data comprising images that each contain up to 5 horizontal lines and randomly-sized circles. Lines and circles were placed at random positions and their pixel intensity was drawn from a uniform distribution between 1 and 10. Examples were generated in an on-line fashion during training. Two illustrative test examples, along with their Fourier-domain measurement representations, are given in Figure 3 (a,b). We compare the results to those obtained using three fixed partial Fourier measurement baselines, following uniform, random, and low-pass subsampling patterns, respectively.

Task model Image reconstruction from partial Fourier measurements was performed by following the methodology in Zhu et al. (2018). We use a deep neural network consisting of two subsequent fully-connected layers with tanh activations that map the 2N (zero-filled) Fourier coefficients (stacked real and imaginary values) to a vector of N pixel values. This vector was subsequently reshaped into a $\sqrt{N} \times \sqrt{N}$ image, and processed by two convolutional layers (ReLU activations, 5×5 kernels, 64 channels, followed by a final convolutional layer (linear activation, 7×7 kernel) that maps these channels to a single $\sqrt{N} \times \sqrt{N}$ image.

Training details The optimization problem was the same as in Equation 11, however with the negative log-likelihood cost \mathcal{L}_s defined as:

$$\mathcal{L}_s = \mathbb{E}_{(\mathbf{x}, \mathbf{s}) \sim p_{\mathcal{D}}} \| f_{\theta}(\mathbf{A}_{(\Phi)}\mathbf{x}) - \mathbf{s} \|_2^2, \tag{12}$$

where we make the parameterization of $\bf A$ by $\bf \Phi$ explicit. The learning rates for $\bf \Phi$ and θ were 1e-3 and 1e-4, respectively. Training was performed for 200,000 iterations across batches of 128 newly generated samples. All other hyperparameters were the same as in Sec. 4.1

Results An overview of the results is given in fig. 3. As expected, uniform subsampling leads to strong spatial aliasing that can not be recovered by the task-model, violating the Nyquist criterion. In turn, random subsampling introduces an incoherent aliasing pattern, that can only in part be recovered. Although not suffering from aliasing, low-pass sampling deteriorates resolution, of which the effect is particularly evident for the broadband/sharp horizontal lines. In contrast, task-adaptive sampling yields high-resolution accurate reconstructions with an improved PSNR (30.4 dB compared to 16.7 dB, 27.4 dB, and 28.7 dB for uniform, random and low-pass sampling, respectively). Note how the sampling pattern (top row) has evolved from a random initialization, similar to that of (d), to the ultimate task-adaptive scheme (f).

4.3 CIFAR10 IMAGE RECONSTRUCTION

Experiment setup The CIFAR10 database (Krizhevsky et al., 2009) contains 60,000 images of 32×32 pixels in 10 different classes. We converted all images to grayscale, and subsequently split them into 50,000 training images, 5,000 validation and 5,000 test images. We again learn partial Fourier sampling and image reconstruction, keeping only 12.5% of the Fourier coefficients, and compare ourselves to the three fixed partial Fourier measurement baselines described in Sec. 4.2.

Task model The challenging reconstruction task for CIFAR10 motivates the adoption of a structured model to enable strong reconstruction. We draw inspiration from iterative proximal-gradient schemes (Parikh et al., 2014) which are dedicated to solving the ill-posed linear measurement problem in equation 1. To that end, we unfold K=5 such iterations, learning an adequate image-domain proximal mapping $\mathcal{P}_{\theta}^{(k)}$ and stepsize $\alpha^{(k)}$ at each fold:

$$\hat{\mathbf{s}}^{(k+1)} = \mathcal{P}_{\theta}^{(k)} \left(\hat{\mathbf{s}}^{(k)} - \alpha^{(k)} \mathbf{F}^* \mathbf{A}_{(\mathbf{\Phi})}^T \left(\mathbf{A}_{(\mathbf{\Phi})} \mathbf{F} \hat{\mathbf{s}}^{(k)} - \mathbf{A}_{(\mathbf{\Phi})} \mathbf{x} \right) \right)$$
(13)

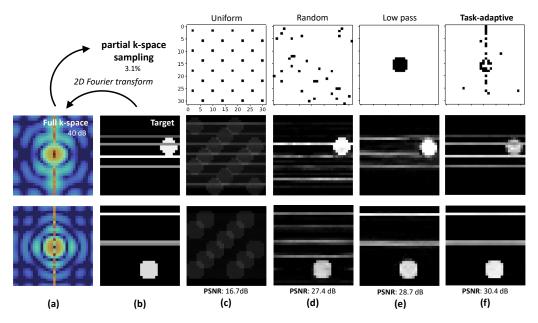


Figure 3: Image reconstruction performance from partial k-space (Fourier) measurements on a custom toy dataset consisting of lines and circles with random locations and sizes. Illustrative examples of the k-space and target images are given in (a,b). The sampling patterns, images and statistical quality metrics for uniform, random, low-pass, and learned task-adaptive sampling are displayed in (c), (d), (e), (f), respectively. In all cases, only 3.1% of the Fourier coefficients have been selected.

where $\mathbf{F} \in \mathbb{R}^{N \times N}$ is a discrete Fourier transform (DFT) matrix, and $(\cdot)^H$ denotes the Hermitian (conjugate transpose). In the above formulation, at each fold a step is taken towards the sampling-consistent subspace that adequately represents the physical measurement of \mathbf{s} by $\mathbf{A}_{(\Phi)}\mathbf{F}$. The trained proximal operator $\mathcal{P}_{\theta}^{(k)}$, a 3-layer convolutional network (3 × 3 kernels, 64 output channels) with ReLU activations followed by a single-output-channel linear 3 × 3 convolutional layer, then projects this onto the manifold of visually plausible images (Mardani et al., 2018), removing noise, aliasing, or blurring artifacts.

Training details Optimization settings were similar to those in Sec. 4.2, leveraging a mean-squared-error (negative log-likelihood) reconstruction cost and a distribution entropy penalty on π_m . To promote visually plausible reconstructions, we added an adversarial (Ledig et al., 2017) cost by adopting a discriminator network $D_{\psi}(\mathbf{s})$ that aims to discriminate between images reconstructed from partial Fourier measurements $\hat{\mathbf{s}}$ and actual images \mathbf{s} . The discriminator comprised 3 convolutional layers (3 × 3 kernels, stride 2, 128 channels) with leaky ReLU activations, followed by global average pooling, 40% dropout, and a logistic binary classification model. The sampling-and task-model parameters were then trained to both minimize the negative log-likelihood cost and maximize the discriminator binary cross-entropy cost \mathcal{L}_D , in addition to the entropy penalty:

$$\hat{\mathbf{\Phi}}, \hat{\theta} = \underset{\mathbf{\Phi}, \theta}{\operatorname{argmin}} \left\{ \mathbb{E}_{(\mathbf{x}, \mathbf{s}) \sim p_{D}} \left\| f_{\theta}(\mathbf{A}_{(\mathbf{\Phi})} \mathbf{x}) - \mathbf{s} \right\|_{2}^{2} + \mu \mathcal{L}_{e} - \lambda \mathcal{L}_{D} \right\}, \tag{14}$$

where λ was set to 0.004. The parameters of the discriminator network, ψ , were jointly optimized with eq. (14) to minimize \mathcal{L}_D . The learning rates for Φ and θ were 1e-3 and 1e-4, respectively. Training was performed for 100 epochs. All other hyperparameters were set as in Sec. 4.1.

Results Figure 4 shows how task-adaptive sampling significantly outperforms all baselines. Both the uniform and random subsampling schemes suffer from severe aliasing which the task-model is not able to adequately restore. While low-pass sampling does not lead to aliasing, the absence of real high-frequency information causes the task-model to 'hallucinate' such (super-resolution) content. This is particularly notable for the example displayed in the bottom row of fig. 4, being less structured (or predictable) than that in the row above it. Instead, task-adaptive sampling learns to select what appear to be pseudo-random samples with a variable spectral density,

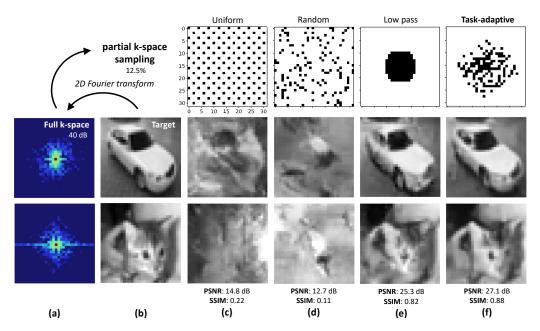


Figure 4: Image reconstruction performance from partial k-space (Fourier) measurements on the CIFAR10 database. Illustrative examples of the k-space and target images are given in (a,b). The sampling patterns, images and statistical quality metrics for uniform, random, low-pass, and learned task-adaptive sampling are displayed in (c), (d), (e), (f), respectively. In all cases, 12.5% of the Fourier coefficients have been selected. Task-adaptive sampling notably outperforms the other approaches, both qualitatively and quantitatively.

thereby acquiring both low and high-frequency information. This allows the task model to produce high-quality, high-resolution reconstructions that go beyond those obtained with the other methods, reaching a PSNR of 27.1 dB and a structural similarity index (SSIM) of 0.88 across the 5000 test examples. For the most competitive fixed-sampling method, low pass, these statistics were 25.3 dB and 0.82, respectively.

5 Conclusions

We have introduced *DPS*, a framework that enables jointly learning a data- and task-driven sampling pattern, with a subsequent task-performing model. The framework is generic and can be combined with any network architecture that performs the required task using the subsampled set of signal elements. Empirically we find the method to perform strongly on toy datasets and canonical deep learning problems. Moreover, the introduced index-conditional parametrization of the subsampling distribution outperforms the fully-factorized Gumbel-top-K method.

Further work can explore the effectiveness of DPS in real-world problems such as MRI scanning. Like all data-driven optimized methods, a learned sampling scheme is at risk of overfitting to a small training dataset. Although we did not observe this issue in our experiments, careful regularization might be required to ensure that this effect is minimal in such high-risk tasks. The fully-differentiable DPS framework allows for flexibility, and interesting extensions can be explored in future work. Rather than learning a fixed set of parameters for the subsampling distribution, a neural network can be used to predict the parameters instead, conditioned on contextual data or the samples acquired so far. Finally, our method currently requires the desired sampling rate to be predetermined as a hyperparameter. Future work can explore if this rate can be jointly optimized to incorporate optimization of the sub-sampling rate.

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A GRADIENT OF GUMBEL-SOFTMAX SAMPLING

For any m^{th} pair of rows $(\mathbf{a}_m, \boldsymbol{\phi}_m)$, any n^{th} element $a_{m,n}$ in \mathbf{a}_m can be differentiated towards all elements in $\boldsymbol{\phi}_m$ through:

$$\nabla_{\phi_{m}} a_{m,n}$$

$$= \nabla_{\phi_{m}} \mathbb{E}_{\mathbf{e}_{m}} \left[\operatorname{softmax}_{\tau} (\phi_{m} + \mathbf{e}_{m}) \Big|_{n} \right]$$

$$= \mathbb{E}_{\mathbf{e}_{m}} \left[\nabla_{\phi_{m}} \operatorname{softmax}_{\tau} (\phi_{m} + \mathbf{e}_{m}) \Big|_{n} \right]$$

$$= \mathbb{E}_{\mathbf{e}_{m}} \left[\nabla_{\phi_{m}} \frac{\exp\{(\phi_{m} + \mathbf{e}_{m})/\tau\}}{\sum_{i=1}^{N} \exp\{(\phi_{m,i} + e_{m,i})/\tau\}} \Big|_{n} \right]$$
(15)

Gumbel noise vector \mathbf{e}_m can be reparametrized as a function of uniform noise vector $\boldsymbol{\epsilon}_m \sim \mathcal{U}(0,1)$ i.i.d., through:

$$\mathbf{e}_m = -\log(-\log(\boldsymbol{\epsilon}_m)). \tag{16}$$

This allows rewriting eq. (15) into:

$$\nabla_{\boldsymbol{\phi}_{m}} a_{m,n} = \mathbb{E}_{\boldsymbol{\epsilon}_{m}} \left[\nabla_{\boldsymbol{\phi}_{m}} \frac{\exp\{(\boldsymbol{\phi}_{m} - \log(-\log(\boldsymbol{\epsilon}_{m})))/\tau\}}{\sum_{i=1}^{N} \exp\{(\boldsymbol{\phi}_{m,i} - \log(-\log(\boldsymbol{\epsilon}_{m,i})))/\tau\}} \right|_{n} \right]$$

$$= \int_{-\infty}^{\infty} \dots \int_{-\infty}^{\infty} \mathbf{P} \left[\boldsymbol{\epsilon}_{m} = [k_{1}, \dots, k_{N}] \right] \nabla_{\boldsymbol{\phi}_{m}} \frac{\exp\{(\boldsymbol{\phi}_{m} - \log(-\log(\boldsymbol{k})))/\tau\}}{\sum_{i=1}^{N} \exp\{(\boldsymbol{\phi}_{m,i} - \log(-\log(k_{i})))/\tau\}} \right|_{n} dk_{N} \dots dk_{1}$$

$$= \int_{-0}^{1} \dots \int_{0}^{1} \mathbf{P} \left[\boldsymbol{\epsilon}_{m,1} = k_{1} \right] \mathbf{P} \left[\boldsymbol{\epsilon}_{m,2} = k_{2} \right] \dots \mathbf{P} \left[\boldsymbol{\epsilon}_{m,N} = k_{N} \right] \cdot \nabla_{\boldsymbol{\phi}_{m}} \frac{\exp\{(\boldsymbol{\phi}_{m} - \log(-\log(\boldsymbol{k})))/\tau\}}{\sum_{i=1}^{N} \exp\{(\boldsymbol{\phi}_{m,i} - \log(-\log(\boldsymbol{k})))/\tau\}} \right|_{n} dk_{N} \dots dk_{1}$$

$$= \int_{-0}^{1} \dots \int_{0}^{1} 1 \cdot \nabla_{\boldsymbol{\phi}_{m}} \frac{\exp\{(\boldsymbol{\phi}_{m} - \log(-\log(\boldsymbol{k})))/\tau\}}{\sum_{i=1}^{N} \exp\{(\boldsymbol{\phi}_{m,i} - \log(-\log(\boldsymbol{k})))/\tau\}} \right|_{n} dk_{N} \dots dk_{1}.$$

$$(17)$$