Deconvolutional Density Network: Modeling Free-Form Conditional Distributions

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Abstract
Conditional density estimation (CDE) is the task of estimating the probability of an event conditioned on some inputs. A neural network (NN) can also be used to compute the output distribution for continuous-domain, which can be viewed as an extension of regression task. Nevertheless, it is difficult to explicitly approximate a distribution without knowing the information of its general form a priori. In order to fit an arbitrary conditional distribution, discretizing the continuous domain into bins is an effective strategy, as long as we have sufficiently narrow bins and very large data. However, collecting enough data is often hard to reach and falls far short of that ideal in many circumstances, especially in multivariate CDE for the curse of dimensionality. In this paper, we demonstrate the benefits of modeling free-form conditional distributions using a deconvolution-based neural net framework, coping with data deficiency problems in discretization. It has the advantage of being flexible but also takes advantage of the hierarchical smoothness offered by the deconvolution layers. We compare our method to a number of other density-estimation approaches and show that our Deconvolutional Density Network (DDN) outperforms the competing methods on many univariate and multivariate tasks. The code of DDN is available at https://github.com/NBICLAB/DDN

Introduction
In supervised discriminative learning, a conditional model learns the conditional distribution function of the outputs given the inputs, while ignoring the distribution of the inputs. If the outputs are random variables, a learned conditional model can map inputs to a discrete set of labels (e.g. classifying handwritten digits), or to values in a continuous domain (e.g. regressing temperature).

What if the expected conditional output is an explicit whole distribution, but not just a value of a random variable? Both classification and regression can be posed as summary statistics of distribution: the location of the maximum (mode) for classification, and the expectation of the distribution for regression. However these do not tell the whole story, and often the entire distribution is much preferred, especially when there are lots of uncertainties, and

\begin{figure}[h]
\centering
\includegraphics[width=\textwidth]{figure1.png}
\caption{Example Conditional Density Estimation, estimating the conditional density on a 2-dimensional domain. This dataset has 35,000 training samples in total, drawn from a distribution constructed based on an animation, by taking each normalized frame and time as conditional density function and condition, respectively. Each sample point consists of a single input $x$ and two outputs $(y^{(1)}, y^{(2)})$. We wish to estimate the conditional density over all possible values of $(y^{(1)}, y^{(2)})$, given different input/condition $x$. Given different value of input, the sub-figures at top and at bottom represent the scatter plot of training samples and the estimated conditional densities, respectively. Notice that the density can change markedly, depending on the input. The density is estimated by DDN.}
\end{figure}

the distributions are broad or multi-modal. For instance, we can predict the distribution curve over a student’s all possible scores given his personality and previous scores, or to estimate the distribution surface over all drop-off locations given pickup location and the time of day (Dutordoir et al. 2018).

This is related to the Conditional Density Estimation (CDE), a more general discriminative learning task, where the goal is to return the distribution over a range of output values. More precisely, given an input $x$, estimate $p(y|x)$, where $y \in \mathcal{D}$ and $\mathcal{D}$ stands for the domain of $y$. Fig. 1 illustrates an example of CDE.

We focus on explicit CDE for fast and direct access to the likelihoods and the whole landscape of distribution. A critical problem with learning an explicit CDE model is the

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ground truth distribution as training signal is often absent. Instead, we can maximize the likelihood. For example, if $y$ belongs to discrete-domain, one could take a neural network classifier using softmax as output layer and cross-entropy as loss function to learn an explicit estimation.

However, if $y$ belongs to continuous-domain, constructing a distribution for explicit CDE tasks becomes not that straightforward. In theory, an infinite number of control parameters is needed to specify arbitrary probability density functions (PDFs) at continuous-domain. This is not practically possible, and is further exacerbated by the finite size of datasets used to tune the parameters.

The only option is to limit ourselves to a finite-dimensional subspace for our output functions. For example, the model could learn to output the mean and standard deviation of a Gaussian distribution. Or, we could output a number of such combinations of parameters, constructing a PDF as a Gaussian mixture model (Bishop 1994) or as normalizing flows (Trippe and Turner 2018).

Hence, we need to apply some assumptions, so-called inductive biases, to the CDE model for continuous-domain. Naturally, it is helpful if one knows the general form of the target distribution a priori. However, there are many cases in which that information is not available, and our chosen form of inductive bias might not fit the underlying distribution very well. Therefore, considering the distribution diversity, we would like to make as few assumptions about the distribution as possible to construct a free-form PDF.

One effective way to construct such free-form PDF curve is piecewise-constant function, in which we specify the values of the function over a set of discrete bins. Ideally, if each bin is narrow enough and has sufficient training samples to draw from, then the law of large numbers should yield a piecewise-constant function that faithfully reflects the underlying true distribution. Despite being a parametric function, this provides extensive flexibility. Moreover, compared to mixture density models, it is easy to learn and produces better performance, as it avoids parts of PDF moving beyond a predefined region (Van Den Oord, Kalchbrenner, and Kavukcuoglu 2016). In fact, the effectiveness of piecewise-constant function for constructing free-form distribution has been verified in generative learning, like PixelRNN (Van Den Oord, Kalchbrenner, and Kavukcuoglu 2016) and Gated PixelCNN (Van Den Oord et al. 2016).

Motivation Unfortunately, the realities of dataset sizes fall far short of that ideal in many circumstances, especially when we only have few data samples. Actually, collecting enough data is hard to reach for many fields in many circumstances. In addition, if we want to finely approximate a conditional density function using many narrow bins (Li, Reich, and Bondell 2021), the data deficiency problem would also lead to poor performance. This problem is even more notable for multivariate CDE on account of the curse of dimensionality, especially when $x$ is in continuous-domain.

Therefore, a question arises here is: can we construct a piecewise-constant function to approximate free-form PDF, on the condition that only few samples are available for learning?

Contribution The answer is yes. The shortage of data in each bin can be mitigated by applying a hierarchical smoothness constraint, which is the main objective of this research.

In this paper, a Deconvolutional Density Network (DDN) is proposed to estimate free-form explicit conditional distributions. It adopts deconvolution (Yu et al. 2017) (including forms of convolution with upsampling (Shelhamer, Long, and Darrell 2017)) to naturally encourage smoothness because of its built-in hierarchical spatial correlation properties. The contribution is decomposed as follows:

1. This work designs a deconvolution-based neural net framework for CDE to explicitly estimate free-form PDF for continuous-domain, especially coping with the data deficiency problem.

2. Experiments manifest DDN achieves the best results on diversified univariate and multivariate CDE tasks in comparison with six different explicit CDE models.

Related Work

This section reviews previous work on CDE and describes the motivations behind our method.

Conceivably kernel-based methods are among the most popular non-neural CDE techniques, as they involve combining kernel functions with data to calculate the probability distributions by interpolating between observed points to predict the probability density of the unseen points. One of the well-recognized kernel-based approaches to estimate CDE is to use the kernel density estimator (KDE) (tsybakov 2008). In principle, KDE can approximate any probability distribution without any prior assumptions. The kernel quantile regression (KQR) method can predict the percentiles of the conditional distribution (Takeuchi et al. 2006; Li, Liu, and Zhu 2007; Takeuchi, Nomura, and Kanamori 2009). A straightforward CDE method, least-squares conditional density estimation (LSCDE), was proposed by Sugiyama et al., employing the squared loss function.

Neural networks have been used for CDE as a result of their expressive power. Mixture density networks (MDNs) (Bishop 1994) are among the earliest strategies for modeling conditional densities as a mixture of distributions, employing NNs to estimate the parameters of each distribution (e.g., mean, standard deviation, etc.). To solve the quadratically growing problem of parameter number in MDNs, the real-valued neural autoregressive density-estimator (RNADE) (Uria, Murray, and Larochelle 2013) applies one-dimensional MDNs to unsupervised density estimation in a practical way.

Another big family of neural density estimators that are flexible is normalizing flows (Tabak and Vanden-Eijnden 2010; Tabak and Turner 2013; Rezende and Mohamed 2015), involving applying different invertible transformations to a simple density until a desired level of complexity is reached. They provide the PDF with a tractable form. In terms of estimating conditional density, Trippe and Turner employ normalising flows as a flexible likelihood model to fit complex conditional densities (Trippe and Turner 2018).

A model agnostic normalizing flow based noise regularization method to estimate conditional density was proposed...
by [Rothfuss et al.] which controls random perturbations to the data during training [Rothfuss et al. 2020]. Utilizing normalizing flow to model conditional distributions, the inverse problems employing maximization of the posterior likelihood is investigated [Xiao, Yan, and Amit 2019]. The Conditional Normalizing Flows (CNFs), where the base models to output space mapping is conditioned on an input, is presented by [Winkler et al.] to model conditional densities. A normalizing flow based autoregressive model (Papamakarios, Pavlakou, and Murray 2017) was proposed for density estimation, by stacking multiple models with similar characteristics.

The study of latent space models has also generated considerable interest in CDE over the past few years. The Gaussian Process CDE (Dutoit et al. 2018) uses latent variables to extend input and maps augmented input onto samples from the conditional distribution by Gaussian processes. As a variant of variational autoencoder (Kingma and Welling 2019), the conditional variational autoencoder (CVAE) is another way to approach CDE (Sohn, Lee, and Yan 2015).

There are also other miscellaneous neural network techniques for CDE. The kernel mixture network (KMN) (Ambrogioni et al. 2017) combines non-parametric and parametric elements by formulating the CDE model as a linear combination of a family of kernels. On the basis of score matching, Sasaki and Hyvärinen designed a neural-kernelized conditional density estimator (NKC) to estimate the conditional density. A conditional sum-product networks (CSPNs) is proposed (Shao et al. 2020) for multivariate and potentially hybrid domains that allow harnessing the expressive power of neural networks while still maintaining tractability guarantees. A neural noise-regularization and data-normalization scheme for CDE was used for financial applications, addressing problems like overfitting, weight initialization, and hyperparameter sensitivity (Rothfuss et al. 2019). Strauss and Oliva devised an energy-based approach for conditional density estimation to eliminate restrictive bias from standard alternatives (Strauss and Oliva 2021).

If we move our focus to generative learning temporally, there are some flexible approaches transforming the continuous-domain density estimation tasks to a discrete domain by piecewise-constant function. The PixelRNNs [Van Den Oord, Kalchbrenner, and Kavukcuoglu 2016] sequentially predicts the pixels in an image by treating the pixel values as discrete random variables using a softmax layer in the conditional distributions. A PixelCNN is also proposed by Van Den Oord, Kalchbrenner, and Kavukcuoglu experimentally found the piecewise-constant function-based CDE is “easy to learn and to produce better performance compared to a continuous distribution” [Van Den Oord, Kalchbrenner, and Kavukcuoglu 2016].

But the strain between infinite-dimensional distributions and finite datasets makes it difficult to faithfully estimate a distribution by piecewise-constant function, especially in scenarios when we only have a small number of samples. This would produce distributions that can be very nonsmooth. Furthermore, the absence of hierarchical smoothness in these methods may also result in failure when predicting complex distribution on a small training set. In addition, the PixelRNN [Van Den Oord, Kalchbrenner, and Kavukcuoglu 2016] and JBCENN [Li, Reich, and Bondell 2021] do not investigate on the exponential computational cost problem when using a piecewise-constant function for multivariate tasks.

In DDN, assuming the true PDFs are commonly smooth and have few spikes, we can use multi-layer deconvolution to hierarchically construct distributions that exhibit spatial coherence at multiple scales, coping with data deficiency problems. Furthermore, a conditional masking-based univariate distribution estimation strategy is adopted to construct an auto-regressive multivariate CDE model.

Methodology

The complete framework of the proposed Deconvolutional Density Network (DDN) is shown in Fig. 2 combining five
successive modules into a single framework. This framework enables DDN to flexibly and explicitly estimate conditional density functions for continuous-domain, reducing needs to large size data. We will describe how the DDN works from Fig. 2(e) back to Fig. 2(a), subsection by subsection.

**Baseline Model**

Here we briefly describe the baseline model, and other components of DDN will be added into it progressively.

If we are given a set of \( M \) input/output samples \( (x_m, y_m) \), then the likelihood of observing those samples under our model is,

\[
\prod_{m=1}^{M} g(y_m; f(x_m, \theta)) .
\]

Taking the log of that likelihood, and negating it, yields the negative log-likelihood,

\[- \sum_{m=1}^{M} \log g(y_m; f(x_m, \theta)) .
\]

Our NN model can be trained by attempting to minimize the expected negative log-likelihood of the training data,

\[
\arg \min_{\theta} E_{(x_m, y_m)} \left[ - \log g(y_m; f(x_m, \theta)) \right] = \arg \min_{\theta} \left[ - \frac{1}{M} \sum_{m=1}^{M} \log g(y_m; f(x_m, \theta)) \right] . \tag{1}
\]

**Partitioning**

Partitioning is the first step for discretizing PDF of \( y \) of continuous-domain (see Fig. 2(e)). In this subsection, we start from the univariate task, i.e., \( y \) is one dimensional. We assume \( \mathcal{D} \), the domain of the variable \( y \), is finite and partitioned into a number of uniform bins. These bins will be used to construct the CDE as a piecewise-constant function. Discretizing \( y \), we partition the output space into \( N \) bins, and denote the \( i \)-th bin using \( B_i \), with uniform bin width \( \Delta B \). If our neural network outputs \( N \) values, \( f_i \) for \( i = 1, \ldots, N \), then we simply interpret \( f_i(x, \theta) \) as the estimated probability that \( y \) falls in the \( i \)-th bin. Then, the estimated conditional density function can be expressed as a piecewise-constant function,

\[
g(y \mid x) = \begin{cases} f_i(x, \theta) / \Delta B & \text{if } y \in B_i \\ \vdots & \text{if } y \in B_i \end{cases} .
\]

We can write the negative log-likelihood loss function as,

\[
\arg \min_{\theta} \left[ - \frac{1}{M} \sum_{m=1}^{M} \log g(y_m \mid x_m) \right] \tag{2}
\]

\[
= \arg \min_{\theta} \left[ - \frac{1}{M} \sum_{m=1}^{M} \sum_{i=1}^{N} \log (f_i(x_m, \theta)) \mathbf{1}_{B_i}(y_m) \right] . \tag{3}
\]

where \( \mathbf{1}_{B_i}(y) \) is an indicator function, equal to 1 when \( y \in B_i \), and zero otherwise.

Let’s move to the multivariate CDE tasks now, i.e., \( y \) is \( J \)-dimensional. In DDN, we translate a \( J \)-dimensional CDE task into \( J \) one-dimensional CDE tasks at first, and take the chain-rule of probability.
and Oliva [2021] to compose them as
\[
p(y^{(1)}, \ldots, y^{(J)}|x) = \prod_{j=1}^{J} p(y^{(j)}|x, y^{(1)}, \ldots, y^{(j-1)}) .
\]

Suppose the network has \( J \) different sets of output nodes, where the number of nodes in each set is \( N \), corresponding to \( y^{(j)} \), the \( j \)\(^{th} \) variable of \( y \). Fig. 2 shows three such sets of outputs. Then we can include all of these outputs in our loss function from Eq. (3) by averaging over them all. Thus, our loss function for multivariate tasks includes a new average over \( J \),

\[
\arg\min_{\theta} \left[ -\frac{1}{M} \sum_{m=1}^{M} \sum_{j=1}^{J} \sum_{i=1}^{N} \log (f_{ij}(x_m, \theta)) \mathbf{1}_{B_{ij}}(y_m^{(j)}) \right],
\]

where \( f_{ij}(x, \theta) \) is the output of the \( i \)\(^{th} \) node in the \( j \)\(^{th} \) set of output nodes, and \( \mathbf{1}_{B_{ij}}(y^{(j)}) \) is the indicator variable that equals 1 when \( y^{(j)} \in B_{ij} \), and zero otherwise.

**Deconvolutional Estimator**

In this subsection, we focus on Fig. 2(d), which is the critical part of our deconvolutional density network. Despite its potential to approximate a free-form conditional distribution, discretizing the continuous-domain distribution by a piecewise-constant function often suffers from statistical undersampling in its bins, resulting in a distribution that is very spiky. Moreover, the conditional density in real-world can be multi-modal and very complex.

If we can assume that the true PDFs are commonly smooth and have few spikes, we could include a hierarchical mechanism that correlates nearby outputs. In this case, hierarchical deconvolution is a strong candidate for modeling complex landscape, as it can introduce smoothness and structural details layer by layer.

In this work, a multi-layer deconvolutional estimator \( h \) is designed to instill spatial coherence at different scales. This has the advantage of constructing the PDF in a multiscaled, hierarchical manner through multiple deconvolution layers, allowing the model to introduce structure and smoothness at different scales.

Note that our deconvolutional layers are not there to invert or undo a convolutional layer. Instead, the deconvolutional layers [Yu et al. 2017] (also known as transposed convolution or upsampling convolution (Shelhamer, Long, and Darrell 2017)) are used to build the output layer-by-layer, in progressively higher resolution.

In our deconvolutional estimator \( h \), the latent vector \( z \) (see next subsection) is firstly mapped to a fully connected layer, which is further reshaped into a multi-channel feature map. These feature maps are split into \( J \) groups, each of which corresponds to a variable of \( y \). Note that, each variable of \( y \) has an independent pathway from the reshaped initial feature map to its own distribution.

The reshaped initial feature maps are then fed into a sequence of deconvolution layers, constructed by upsampling and convolution. Upsampling expands the feature map by an integer scale factor, while the convolution transforms the expanded feature map with shared weights introducing spatial correlation. After these deconvolution layers, we obtain a vector of unnormalized logits, which are finally fed into a SoftMax layer to get a discretized probability vectors. At this point, we have \( J \) discretized probability vectors, which are used in the loss functions (5).

**Encoder-Estimator Framework**

Let’s move our attention to Fig. 2(b & c) now. Although data insufficiency problem in piecewise-constant function can be mitigated by deconvolution. There are still some remaining questions.

What if we only have extremely small number of samples for learning? What if the target PDF is not that smooth? How to calibrate the estimated distribution?

We wish the DDN to focus on the statistically salient features of limited data, not the contingencies specific to our particular training samples.

Since the layer by layer abstraction enables neural network to extract underlying features of data, adding noise to each sample at latent space could provide the model with lots of unseen but reasonable patterns during learning, to solve the data deficiency problem.

Therefore, we split our network into two main parts: an encoder \( q(z|x) \) (see Fig. 2(b)), and an estimator \( h(y|z) \); a latent space (see Fig. 2(c)) sits between the two.

This is inspired by the Deep VIB framework [Alemi et al. 2017; Alemi, Fischer, and Dillon 2018]. The encoder outputs two equal-sized vectors, \( \mu \) and \( \sigma \), representing the mean and standard deviation for a multivariate, diagonal Gaussian distribution. The latent vectors \( z \) are stochastic, drawn from that Gaussian distribution using the reparameterization trick, \( z = \mu + \sigma \cdot \epsilon \), where \( \epsilon \) is drawn from the standard Normal distribution, \( \epsilon \sim \mathcal{N}(0, I) \). In order to seek a less complex and more economical latent representation, we could force the network to sacrifice some encoding fidelity [Alemi et al. 2017] by adding KL-divergence term \( L_{KL} \) to loss function, where \( L_{KL} = D_{KL}[q(z|x) \| \mathcal{N}(0, 1)] \) and \( q(z|x) \sim \mathcal{N}(\mu(x; \theta), \sigma^2(x; \theta)) \). In DDN, we call this latent-space layer a variational layer (VL).

The adoption of Encoder-Estimator framework makes the final objective function of DDN be written as

\[
\frac{1}{M} \sum_{m=1}^{M} \left\{ E_{z \sim \mathcal{N}(0, 1)} \left[ -\frac{1}{J} \sum_{j=1}^{J} \sum_{i=1}^{N} \log (h_{ij}(y^{(j)}_{m}|z)) \right] \times \mathbf{1}_{B_{ij}}(y^{(j)}_{m}) + \beta L_{KL} \right\}.
\]

Recall that \( h_{ij} \) and \( q \) are both functions of the network parameters, \( \theta \). The \( \beta \) coefficient influences on the calibration of estimated distribution of DDN directly [Alemi et al. 2017; Higgins et al. 2017]. When \( \beta = 0 \), the latent-space distribution can relax, and spread out, allowing the samples to disambiguate from each other. Increasing \( \beta \) encourages the latent-space distribution to contract toward the origin [Alemi et al. 2017]. In this way, \( \beta \) influences the calibrating estimated distribution [Alemi, Fischer, and Dillon 2018] by forcing the model to focus on the statistically salient features, and thereby help to improve the confidence.
Notice that, different from CVAE (Sohn, Lee, and Yan [2015]) or other related latent space CDE models, the DDN doesn’t need Monte Carlo sampling on \( z \) to access \( p(y|x) \). In addition, one may argue that we may have to integrate over \( z \) to obtain an explicit expression for \( p(y|x) \), which is generally intractable. Instead, during testing of DDN, we simply pass the mean vector to the deconvolutional estimator \( q \). In practice, this is realized by setting \( \epsilon \) to zero to remove latent-space stochasticity.

**Conditional Masking**

This subsection is about Fig. 2(a). The DDN relies on partitioning the domain of \( y \) into uniform bins, which has a computational cost exponential in \( J \), where \( J \) is the dimension of \( y \). For multivariate tasks, in order to avoid the number of bins produce scales exponentially with the dimension, we adopt auto-regressive structure (Larochelle and Murray [2011] using the chain-rule of probability (recall subsection Partitioning). Given this, how to calculate the univariate distributions in chain-rule of probability?

In this work, we adopt a conditional masking mechanism to configure DDN for predicting different univariate conditional densities, which can be further assembled for the target multivariate density. At the beginning of learning, we first generate \( K \) permutations for \( (1, 2, \ldots, J) \). Given the \( k \)th permutation is defined as \((A_1^k, A_2^k, \ldots, A_J^k)\), according to Eq. (4), the \( k \)th path to calculate the multivariate conditional distribution is

\[
p\left( y^{(1)}, y^{(2)}, \ldots, y^{(J)} | x \right) \\
= \prod_{j=1}^{J} p\left( y^{(A_j^k)} | x, y^{(A_1^k)}, \ldots, y^{(A_{j-1}^k)} \right). \tag{7}
\]

We further define a conditional mask vector set \( S \) for masking \( y \), each element of which is a \( J \)-dimensional binary vector. Each mask vector corresponds to a condition in a path, i.e., \((y^{(A_1^k)}, \ldots, y^{(A_{J-1}^k)})\) except \( x \). Note the duplicate conditions in \( S \) should be reduced. For a condition and its corresponding mask vector, the \( j \)th component of mask is set to one if \( y^{(j)} \) belongs to this condition, otherwise it is set to zero.

During learning, given a sample \((x_m, y_m)\), we randomly choose a mask \( rm \) from \( S \) at first, and concatenate \( x_m \), \( rm \), and \( y_m \) as the input to DDN. In this way, we can easily assemble different univariate conditional models into a single DDN framework. For example, a conditional binary mask \((1, 0, 1)\) for \( y_m \) means we want to set \((x_m, y_m^{(1)}, y_m^{(3)})\) as a condition.

The mask actually plays a role of information gate, i.e., one means letting it in and zero means keeping it out. Particularly, when \( y_m \cdot rm \) is zero, mask \( rm \) is one means the zero comes from domain. Otherwise, if mask \( rm \) is zero, the information of \( y_m \) is blocked.

In test, we continue to take this conditional masking mechanism on DDN to access different univariate distributions and send them into Eq. (7) to calculate the multivariate conditional distribution. In addition, we freeze chain-rule paths, i.e. the mask vector set \( S \), at the inference to ensure a normalized distribution. Ideally, the predicted distribution by different paths should be the same as each other. However, in practice, we take the average joint distribution from all \( K \) paths to reduce predicting errors.

**Experiments**

This section examines the efficacy of DDN for explicitly estimating free-form conditional PDF on continuous-domain. The following questions are intended to be addressed through our experiments: (a) What is the influence of each of the main components of DDN on the model’s performance? (b) What is the influence of the key hyperparameter \( \beta \) on the estimated density? (c) Is light of the ground truth distributions, how does DDN compare to competitive techniques on low-dimensional toy tasks? (d) In real-world tasks, what is DDN’s performance compared to the competitors?

In response to these questions, we first designed ablation experiments to demonstrate the need for different components of the proposed DDN, followed by an analysis of the hyperparameter \( \beta \). Then we performed experiments on four toy tasks with known ground truth conditional distributions \( p(y|x) \) and evaluated the model using several univariate and multivariate benchmark datasets from real-world applications.

**Experimental Configuration**

**Toy & Real World Tasks**

Considering the ground truth conditional distributions are important for evaluating the accuracy, we designed four 2D toy tasks which \( p(y|x) \) can be easily calculated (See Fig. 5), where \( x \sim U(-1, 1) \). The details of generating training samples are as follows:

- **Squares**: For given \( x \); \( \lambda \sim \text{Bern}(0.5); \quad a^{(1)}, a^{(2)} \sim U(-5 + x, -1 + x); b^{(1)}, b^{(2)} \sim U(1 - x, 5 - x); y^{(1)} = \lambda \ast a^{(1)} + (1 - \lambda) \ast b^{(1)}; y^{(2)} = \lambda \ast a^{(2)} + (1 - \lambda) \ast b^{(2)} \)
- **Half Gaussian**: For given \( x \); \( a, b \sim N(0, 2); y^{(1)} = |a| \cos \pi x - b \sin \pi x; y^{(2)} = |a| \sin \pi x + b \cos \pi x \)
- **Gaussian Stick**: For given \( x \); \( a \sim N(0, 1); b \sim U(-6, 6); c = (-0.75 + x) / 2; y^{(1)} = a \cos c \pi - b \sin c \pi; y^{(2)} = a \sin c \pi + b \cos c \pi \)
- **Elastic Ring**: For given \( x \); \( d \sim U(0, 2); \theta \sim U(0, 2 \pi); y^{(1)} = (4 + 2 \ast x + d) \cos \theta; y^{(2)} = (4 - 2 \ast x + d) \ast \sin \theta \)

The validation also included seven real-world datasets from the UCI machine learning repository [Dua and Graff 2017]. We standardized all inputs and targets using z-score normalization. Table 1 shows the characteristics of each dataset. All experiments were running on a Tesla V100 GPU. We used Adam optimizer with the default configuration. The learning rate was set to 3e-4, and the batch size was 256.
Models In all experiments, the estimator of our DDN contained two Upsample-Conv-BatchNorm-LeakyReLU blocks and one Upsample-Conv block in each target dimension. We configured the blocks so that the output space was partitioned into \( N = 256 \) bins. We set the number of permutations \( K \) to the maximum number of permutations \( K_{\text{max}} \) if \( K_{\text{max}} \) does not exceed 5; otherwise, \( K \) was set to 5. We used 16 latent codes in VL. For 1D tasks, DDN is not using the masking strategy, which means that its input only condition \( x \).

We conducted experiments to compare DDN with seven other neural-network-based CDE methods, including DDN without VL, MDN (Bishop 1994) employing 20 Gaussian kernels, MAF (Papamakarios, Pavlakou, and Murray 2017) using 5 autoregressive transformations, NSF (Durkan et al. 2019) adopting autoregressive model using 5 RQS transformations, RNF utilizing 5 radial transformations (Trippe and Turner 2018), JBCENN (Li, Reich, and Bondell 2021) using 256 bins similar to our DDN model, and Multilayer Perceptron with Softmax (MLP). For each model, most hyperparameters were determined by recommendation, but the ones which didn’t work well on given tasks were tuned by trial and error.

All models began with the same basic three-layers network (input layer and two hidden layers) with 64 tanh neurons in each hidden layer. We normalized every hidden layer using batch-normalization. To balance the proportion of conditional information and masked information on twodimensional and higher-dimensional tasks for our DDN, we used two independent linear layers with 32 tanh neurons to replace the first layer in the basic network, which map \( x \) and \([y \cdot m, m]\) to vectors with same length 32, where \( x \) is the condition, \( y \) is the target value, \( m \) is the mask used in DDN, and \([::]\) denotes the concatenation operation.

It is necessary to set the range of target variables when using DDN, JBCENN, and MLP. We set the range to \((-10, 10)\) for the toy tasks. For real-world tasks, the range was set to \([\lfloor y_{\text{min}} \rfloor - 1, \lceil y_{\text{max}} \rceil + 1]\), where \( y_{\text{min}} \) and \( y_{\text{max}} \) denote the minimum and maximum values of the target random variable \( y \), respectively. We set the probability density to 0 for \( y \)-values outside of that range for these three methods.

MAF in this paper was constructed as described in (Papamakarios, Pavlakou, and Murray 2017). The basic network with an additional linear layer was used to parameterize the transformation function \( z^{(i)} = f(x, y^{(i)}) = (y^{(i)} - \mu(x, y^{(i-1)}))/\exp(\alpha(x, y^{(i-1)})), \) where \( x \) is the condition and \( y \) is the target. Note that for one N-dimensional CDE task, a MAF was parameterized by \( N \) basic networks.

NSF was constructed using the same framework as the MAF, whose transformation function was \( z^{(i)} = f(x, y^{(i)}) = RQS_{\theta}(x, y^{(i-1)})(y^{(i)}), \) where RQS is the rational-quadratic spline function from (Durkan et al. 2019). For all toy tasks, we set the interval to \((-10, 10)\) for the RQS function; for all UCI tasks, we set it to \((-3, 3)\).

The RNF framework was the same as MAF, but its transformation function was replaced by the radial function in (Trippe and Turner 2018).

Evaluation Metrics For real-world tasks, we used the log-likelihood (LL) to evaluate a CDE model. Moreover, we reported the standard deviation of the log-likelihood to demonstrate the stability of DDN in comparison to other techniques. Nevertheless, for the toy tasks, it’s feasible to directly measure the difference between the estimated distribution and the ground truth, as the ground truth distribution \( p(y|x) \) can be easily derived. We used the sum of squared errors (SSE) for measuring the difference.

Ablation Experiment We designed the ablation experiments on Elastic Ring task by comparing DDN, DDN without the variational layer, MLP, MLP with variational layer, and MLP with 10x data samples in estimating free-form conditional distributions. The DDN without VL was designed to show the influence of VL, as we replaced VL with an activated dense layer. Moreover, the MLP, estimating a distribution by piecewise-constant function but without deconvolution, was used to show the influence of the deconvolutional estimator. The estimated densities when \( x = 0.75 \) are shown in Fig. 3.
Ground Truth \hspace{1cm} \beta = 0.1 \hspace{1cm} \beta = 0.05 \\
\begin{array}{cccc}
\beta = 0.01 \hspace{1cm} \beta = 0.005 \hspace{1cm} \text{no VL}
\end{array}

Figure 4: Comparison of estimated densities of DDNs with different $\beta$ values on Gaussian Stick task, where the condition $x = -0.75$.

are provided in each bin. Although coupling with VL or training on 10x samples can mitigate this problem, the estimated distribution is still spiky. The problem can be exacerbated with the increase of dimension due to the curse of dimensionality. VL exposes MLP at more unseen patterns in latent space during learning, but it cannot guarantee the smoothness of density. By contrast, the DDN mitigates this problem by employing deconvolution to introduce spatial correlation hierarchically. Furthermore, note that VL alleviates the over-fitting problem of DDN, and enables calibration. Thus, both deconvolution and VL contribute something to the performance of the overall system.

**Analysis for Hyperparameter $\beta$**

We also studied the impact of different $\beta$ on estimated distribution. As mentioned in Alemi et al. (2017), $\beta$ is used to control the objective of information bottleneck, resulting in a compressed representation in latent space, which can help regularize DDN and avoid it overfitting the training data. Fig. 4 illustrates the larger $\beta$ can diffuse the estimated density and smoothen the function. However, an over-large $\beta$ would degrade the fidelity of an estimated density, as it puts too much weight on the homogeneity of latent variables, breaking the inter-correlation between condition and estimated density. Notice that if the available data is very sufficient, a smaller $\beta$ is suggested to avoid the over-smooth surface of a density function.

**Performance**

**Toy Tasks** For each 2D toy task, we randomly generated 2000 samples from the joint probability distributions $p(x, y^{(1)}, y^{(2)})$ for training. Moreover, we trained all models over 5000 epochs. We calculated the SSE between the ground truth conditional density and the models’ estimated densities.

The results on the 2D toy task are shown in Fig. 5. Moreover, the average SSE over four conditions is provided for each task. It can be observed that the DDN could adapt its estimation to diverse densities, and achieves the best average rank, which is 1.5. The closest competitor is RNF, with a score 2.25. Furthermore, there is only marginal difference the MAF and RNF are ahead of DDN at two tasks. These results verify the shapeshifting ability of DDN.

**Real World Tasks** We also compared DDN with other six CDE techniques on real world datasets. In addition, DDN with different $\beta$ values ($\beta \in \{0.5, 0.1, 0.02\}$) and DDN without VL are all reported for comparison. The models were trained independently for 10 trials on each dataset, with 3000 epochs per experiment. For each trial, we randomly split the data into a training set and test set with a ratio of 3:7, to mimic data deficiency scenario.

Table 2 reports the average log-likelihood and standard deviation of log-likelihood over 10 trials. The DDNs achieve the best average log-likelihood among all the methods. It should be noted that the DDN (no VL) also yields better performance, but VL further improves it. It verifies the deconvolution layers can aid DDN to adapt the real-world distributions. Notice that the DDNs progressively stretch their leading advantage, as the increase of difference between distributions in higher dimensional space. Moreover, DDN also achieves the lowest standard deviation on four datasets out of five. Among these results, MLP reaches the lowest standard deviation on the Energy, but DDN yields the closest results to MLP on this one. This exhibits the DDN is stabler
in estimating free-form conditional distributions.

**Conclusion**

In this work, we design a deconvolutional density network framework for conditional density estimation to explicitly estimate free-form PDF for continuous-domain, especially coping with the data deficiency problem. Assuming the true PDFs are commonly smooth, we use multi-layer deconvolution to hierarchically construct distributions that exhibit spatial coherence at multiple scales, to mitigate the shortage of data. In the experiment, the DDN model clearly showed the best all-around performance on all tasks.

In the future, much more work is required to unpack the features and to improve the performance of the DDN. For example, how does the deconvolution architecture (number of layers, kernel size) affect the performance? How does the dimension of the latent space depend on the dataset? Is there a better alternative to using piecewise-constant functions for creating our free-form distributions? How to accelerate the convergence of DDN, especially for high-dimensional tasks?

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