Towards Biologically Plausible Deep Learning

Yoshua Bengio¹, Dong-Hyun Lee, Jorg Bornschein, Thomas Mesnard and Zhouhan Lin

Montreal Institute for Learning Algorithms, University of Montreal, Montreal, QC, H3C 3J7 ¹CIFAR Senior Fellow

Abstract

Neuroscientists have long criticised deep learning algorithms as incompatible with current knowledge of neurobiology. We explore more biologically plausible versions of deep representation learning, focusing here mostly on unsupervised learning but developing a learning mechanism that could account for supervised, unsupervised and reinforcement learning. The starting point is that the basic learning rule believed to govern synaptic weight updates (Spike-Timing-Dependent Plasticity) arises out of a simple update rule that makes a lot of sense from a machine learning point of view and can be interpreted as gradient descent on some objective function so long as the neuronal dynamics push firing rates towards better values of the objective function (be it supervised, unsupervised, or rewarddriven). The second main idea is that this corresponds to a form of the variational EM algorithm, i.e., with approximate rather than exact posteriors, implemented by neural dynamics. Another contribution of this paper is that the gradients required for updating the hidden states in the above variational interpretation can be estimated using an approximation that only requires propagating activations forward and backward, with pairs of layers learning to form a denoising auto-encoder. Finally, we extend the theory about the probabilistic interpretation of auto-encoders to justify improved sampling schemes based on the generative interpretation of denoising auto-encoders, and we validate all these ideas on generative learning tasks.

1. Introduction

Deep learning and artificial neural networks have taken their inspiration from brains, but mostly for the form of the computation performed (with much of the biology, such as the presence of spikes remaining to be accounted for). However, what is lacking currently is a credible *machine learning interpretation of the learning rules* that seem to

exist in biological neurons that would explain *efficient* joint training of a deep neural network, i.e., accounting for *credit* assignment through a long chain of neural connections. Solving the credit assignment problem therefore means identifying neurons and weights that are responsible for a desired outcome and changing parameters accordingly. Whereas back-propagation offers a machine learning answer, it is not biologically plausible, as discussed in the next paragraph. Finding a biologically plausible machine learning approach for credit assignment in deep networks is the main long-term question to which this paper contributes.

Let us first consider the claim that state-of-the-art deep learning algorithms rely on mechanisms that seem biologically implausible, such as gradient back-propagation, i.e., the mechanism for computing the gradient of an objective function with respect to neural activations and parameters. The following difficulties can be raised regarding the biological plausibility of back-propagation: (1) the back-propagation computation (coming down from the output layer to lower hidden layers) is purely linear, whereas biological neurons interleave linear and non-linear operations, (2) if the feedback paths known to exist in the brain (with their own synapses and maybe their own neurons) were used to propagate credit assignment by backprop, they would need precise knowledge of the derivatives of the non-linearities at the operating point used in the corresponding feedforward computation on the feedforward path¹, (3) similarly, these feedback paths would have to use exact symmetric weights (with the same connectivity, transposed) of the feedforward connections,² (4) real neurons communicate by (possibly stochastic) binary values (spikes), not by clean continuous values, (5) the computation would have to be precisely clocked to alternate between feedforward and back-propagation phases (since the latter needs the former's results), and (6) it is not clear where the output targets would come from. The approach

¹ and with neurons not all being exactly the same, it could be difficult to match the right estimated derivatives

²this is known as the *weight transport* problem (Lillicrap *et al.*, 2014)

proposed in this paper has the ambition to address all these issues, although some question marks as to a possible biological implementations remain, and of course many more detailed elements of the biology that need to be accounted for are not covered here.

Note that back-propagation is used not just for classical supervised learning but also for many unsupervised learning algorithms, including all kinds of auto-encoders: sparse auto-encoders (Ranzato et al., 2007; Goodfellow et al., 2009), denoising auto-encoders (Vincent et al., 2008), contractive auto-encoders (Rifai et al., 2011), and more recently, variational auto-encoders (Kingma and Welling, 2014). Other unsupervised learning algorithms exist which do not rely on back-propagation, such as the various Boltzmann machine learning algorithms (Hinton and Sejnowski, 1986; Smolensky, 1986; Hinton et al., 2006; Salakhutdinov and Hinton, 2009). Boltzmann machines are probably the most biologically plausible learning algorithms for deep architectures that we currently know, but they also face several question marks in this regard, such as the weight transport problem ((3) above) to achieve symmetric weights, and the positive-phase vs negative-phase synchronization question (similar to (5) above).

Our starting point (Sec. 2) proposes an interpretation of the main learning rule observed in biological synapses: Spike-Timing-Dependent Plasticity (STDP). Inspired by earlier ideas (Xie and Seung, 2000; Hinton, 2007), we first show via both an intuitive argument and a simulation that STDP could be seen as stochastic gradient descent if only the neuron was driven by a feedback signal that either increases or decreases the neuron's firing rate in proportion to the gradient of an objective function with respect to the neuron's voltage potential.

In Sec. 3 we present the first machine learning interpretation of STDP that gives rise to efficient credit assignment through multiple layers. We first argue that the above interpretation of STDP suggests that neural dynamics (which creates the above changes in neuronal activations thanks to feedback and lateral connections) correspond to inference towards neural configurations that are more consistent with each other and with the observations (inputs, targets, or rewards). This view is analogous to the interpretation of inference in Boltzmann machines while avoiding the need to obtain representative samples from the stationary distribution of an MCMC. Going beyond Hinton's proposal, it naturally suggests that the training procedure corresponds to a form of variational EM (Neal and Hinton, 1999) (see Sec.3), possibly based on MAP (maximum a posteriori) or MCMC (Markov Chain Monte-Carlo) approximations. In Sec. 4 we show how this mathematical framework suggests a training procedure for a deep directed generative network with many layers of latent variables. However, the above interpretation would still require to compute some gradients. Another contribution (Sec. 6) is to show that one can estimate these gradients via an approximation that only involves ordinary neural computation and no explicit derivatives, following previous (unpublished) work on target propagation (Bengio, 2014; Lee *et al.*, 2014). We introduce a novel justification for difference target propagation (Lee *et al.*, 2014), exploiting the fact that the proposed learning mechanism can be interpreted as training a denoising auto-encoder. As discussed in Sec. 5 these alternative interpretations of the model provide different ways to sample from it, and we found that better samples could be obtained.

2. STDP as Stochastic Gradient Descent

Spike-Timing-Dependent Plasticity or STDP is believed to be the main form of synaptic change in neurons (Markram and Sakmann, 1995; Gerstner et al., 1996) and it relates the expected change in synaptic weights to the timing difference between post-synaptic spikes and pre-synaptic spikes. Although it is the result of experimental observations in biological neurons, its interpretation as part of a learning procedure that could explain learning in deep networks remains unclear. Xie and Seung (2000) nicely showed how STDP can correspond to a differential anti-Hebbian plasticity, i.e., the synaptic change is proportional the product of pre-synaptic activity and the temporal derivative of the post-synaptic activity. The question is how this could make sense from a machine learning point of view. This paper aims at proposing such an interpretation, starting from the general idea introduced by Hinton (2007), anchoring it in a novel machine learning interpretation, and extending it to deep unsupervised generative modeling of the data.

What has been observed in STDP is that the weights change if there is a pre-synaptic spike in the temporal vicinity of a post-synaptic spike: that change is positive if the post-synaptic spike happens just after the pre-synaptic spike, negative if it happens just before. Furthermore, the amount of change decays to zero as the temporal difference between the two spikes increases in magnitude. We are thus interested in this temporal window around a pre-synaptic spike during which we could have a post-synaptic spike, before or after the pre-synaptic spike.

We propose a novel explanation for the STDP curve as a consequence of an actual update equation which makes a lot of sense from a machine learning perspective:

$$\Delta W_{ij} \propto S_i \dot{V}_i,$$
 (1)

where \dot{V}_j indicates the temporal derivative of V_j , S_i indicates the pre-synaptic spike (from neuron i), and V_j indicates the post-synaptic voltage potential (of neuron j).

To see how the above update rule can give rise to STDP, consider the average effect of the rest of the inputs into

neuron j, which induces an average temporal change \dot{V}_j of the post-synaptic potential (which we assume approximately constant throughout the duration of the window around the pre-synaptic firing event), and assume that at time 0 when S_i spikes, V_j is below the firing threshold. Let us call ΔT the temporal difference between the post-synaptic spike and the pre-synaptic spike.

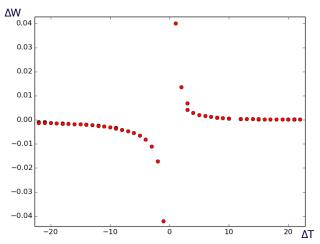


Figure 1. Result of simulation around pre-synaptic spike (time 0) showing indirectly the effect of a change in the rate of change in the post-synaptic voltage, \dot{V}_j on both the average time difference between pre- and post-synaptic spikes (horizontal axis, ΔT) and the average weight change (vertical axis, ΔW_{ij}), when the latter follows Eq. 1. This corresponds very well to the observed relationship between ΔT and ΔW_{ij} in the biological literature.

First, let us consider the case where $V_j > 0$, i.e., the voltage potential increases. Depending on the magnitude of V_j , it will take more or less time for V_j to reach the firing threshold, i.e. more time for smaller V_j . Hence a longer ΔT corresponds to a smaller ΔW_{ij} , and a positive ΔT to a positive ΔW_{ij} , as observed for STDP.

Second, let us consider the case where $V_i < 0$, i.e., the voltage potential has been decreasing (remember that we are considering the average effect of the rest of the inputs into neuron j, and assuming that the temporal slope of that effect is locally constant). Thus, it is likely that earlier before the pre-synaptic spike the post-synaptic voltage V_j had been high enough to be above the firing threshold. How far back in the past again depends monotonically on \dot{V}_j . Hence a negative ΔT corresponds to a negative ΔW_{ij} and a more negative ΔT corresponds to a ΔW_{ij} that is smaller in magnitude. This corresponds perfectly to the kind of relationship that is observed by biologists with STDP (see Figure 7 of Bi and Poo (1998) or Figure 1 of Sjöström and Gerstner (2010), at http://www.scholarpedia.org/article/ Spike-timing_dependent_plasticity). In a simulation inspired by the above analysis, we observe essentially the same curve relating the average ΔT and the ΔW_{ij} that is associated with it, as illustrated in Figure 1.

Clearly, the consequence of Eq. 1 is that if the change ΔV_j corresponds to improving some objective function J, then STDP

corresponds to approximate stochastic gradient descent in that objective function. With this view, STDP would implement the delta rule (gradient descent on a one-layer network) if the post-synaptic activation changes in the direction of the gradient.

3. Variational EM with Learned Approximate Inference

To take advantage of the above statement, the dynamics of the neural network must be such that neural activities move towards better values of some objective function J. Hence we would like to define such an objective function in a way that is consistent with the actual neural computation being performed (for fixed weights W), in the sense that the expected temporal change of the voltage potentials approximately corresponds to increases in J. In this paper, we are going to consider the voltage potentials as the central variables of interest which influence J and consider them as latent variables V (denoted h below to keep machine learning interpretation general), while we will consider the actual spike trains S as non-linear noisy corruptions of V, a form of quantization (with the "noise level" controlled either by the integration time or the number of redundant neurons in an ensemble (Legenstein and Maass, 2014)). This view makes the application of the denoising auto-encoder theorems discussed in Sec. 5 more straightforward.

The main contribution of this paper is to propose and give support to the hypothesis that J comes out of a *variational bound* on the likelihood of the data. Variational bounds have been proposed to justify various learning algorithms for generative models (Hinton *et al.*, 1995) (Sec. 7). To keep the mapping to biology open, consider such bounds and the associated criteria that may be derived from them, using an abstract notation with *observed variable x* and *latent variable h*. If we have a model p(x, h) of their joint distribution, as well as some approximate inference mechanism defining a conditional distribution $q^*(H|x)$, the observed data log-likelihood $\log p(x)$ can be decomposed as

$$\log p(x) = \log p(x) \sum_{h} q^{*}(h|x)$$

$$= \sum_{h} q^{*}(h|x) \log \frac{p(x,h)q^{*}(h|x)}{p(h|x)q^{*}(h|x)}$$

$$= E_{q^{*}(H|x)}[\log p(x,H)] + \mathbf{H}[q^{*}(H|x)]$$

$$+ \mathbf{KL}(q^{*}(H|x)||p(H|x)), \quad (2)$$

where $\mathbf{H}[]$ denotes entropy and $\mathbf{KL}(||)$ the Kullback-Leibler (KL) divergence, and where we have used sums but integrals should be considered when the variables are continuous. Since both the entropy and the KL-divergence are non-negative, we can either bound the log-likelihood via

$$\log p(x) \ge E_{q^*(H|x)}[\log p(x,H)] + \mathbf{H}[q^*(H|x)], \quad (3)$$

or if we care only about optimizing p,

$$\log p(x) \ge E_{q^*(H|x)}[\log p(x, H)]. \tag{4}$$

The idea of variational bounds as proxies for the log-likelihood is that as far as optimizing p is concerned, i.e., dropping the entropy term which does not depend on p, the bound becomes tight when $q^*(H|x) = p(H|x)$. This suggests that $q^*(H|x)$ should approximate p(H|x). Fixing $p^*(H|x) = p(H|x)$ and optimizing p with p fixed is the EM algorithm. Here (and in general) this is not possible so we consider variational methods in which $p^*(H|x)$ approximates but does not reach p(H|x). This variational bound has recently been used to justify another biologically plausible

update rule (Rezende and Gerstner, 2014), which however relied on the REINFORCE algorithm (Williams, 1992) rather than on inference to obtain credit assignment to internal neurons.

We propose to decompose $q^*(H|x)$ in two components: parametric initialization $q_0(H|x) = q(H|x)$ and iterative inference, implicitly defining $q^*(H|x) = q_T(H|x)$ via a deterministic or stochastic update, or transition operator

$$q_t(H|x) = A(x) \ q_{t-1}(H|x).$$
 (5)

The variational bound suggests that A(x) should gradually bring $q_t(H|x)$ closer to p(H|x). At the same time, to make sure that a few steps will be sufficient to approach p(H|x), one may add a term in the objective function to make $q_0(H|x)$ closer to p(H|x), as well as to encourage p(x,h) to favor solutions p(H|x) that can be easily approximated by $q_t(H|x)$ even for small t.

For this purpose, consider as training objective a regularized variational MAP-EM criterion (for a given *x*):

$$J = \log p(x, h) + \alpha \log q(h|x), \tag{6}$$

where h is a free variable (for each x) initialized from q(H|x)and then iteratively updated to approximately maximize J. The total objective function is just the average of J over all examples after having performed inference (the approximate maximization over h for each x). A reasonable variant would not just encourage $q = q_0$ to generate h (given x), but all the q_t 's for t > 0 as well. Alternatively, the iterative inference could be performed by stochastically increasing J, i.e., via a Markov chain which may correspond to probabilistic inference with spiking neurons (Pecevski et al., 2011). The corresponding variational MAP or variational MCMC algorithm would be as in Algorithm 1. For the stochastic version one would inject noise when updating h. Variational MCMC (de Freitas et al., 2001) can be used to approximate the posterior, e.g., as in the model from Salimans et al. (2014). However, a rejection step does not look very biologically plausible (both for the need of returning to a previous state and for the need to evaluate the joint likelihood, a global quantity). On the other hand, a biased MCMC with no rejection step, such as the stochastic gradient Langevin MCMC of Welling and Teh (2011) can work very well in practice.

Algorithm 1 Variational MAP (or MCMC) SGD algorithm for gradually improving the agreement between the values of the latent variables h and the observed data x. q(h|x) is a learned parametric initialization for h, p(h) is a parametric prior on the latent variables, and p(x|h) specifies how to generate x given h. Objective function J is defined in Eq. 6 Learning rates δ and ϵ respectively control the optimization of h and of parameters θ (of both q and p).

Initialize
$$h \sim q(h|x)$$
 for $t=1$ to T do
B $h \leftarrow h + \delta \frac{\partial J}{\partial h}$ (optional: add noise for MCMC) end for $\theta \leftarrow \theta + \epsilon \frac{\partial J}{\partial \theta}$

4. Training a Deep Generative Model

There is strong biological evidence of a distinct pattern of connectivity between cortical areas that distinguishes between "feed-forward" and "feedback" connections (Douglas *et al.*, 1989) at

the level of the microcircuit of cortex (i.e., feedforward and feedback connections do not land in the same type of cells). Furthermore, the feedforward connections form a directed acyclic graph with nodes (areas) updated in a particular order, e.g., in the visual cortex (Felleman and Essen, 1991). So consider Algorithm 1 with h decomposed into multiple layers, with the conditional independence structure of a directed graphical model structured as a chain, both for p (going down) and for q (going up):

$$p(x,h) = p(x|h^{(1)}) \left(\prod_{k=1}^{M-1} p(h^{(k)}|h^{(k+1)}) \right) p(h^{(M)})$$
$$q(h|x) = q(h^{(1)}|x) \prod_{k=1}^{M-1} q(h^{(k+1)}|h^{(k)}). \tag{7}$$

This clearly decouples the updates associated with each layer, for both h and θ , making these updates "local" to the layer k, based on "feedback" from layer k-1 and k+1. Nonetheless, thanks to the iterative nature of the updates of h, all the layers are interacting via both feedforward $(q(h^{(k)}|h^{(k-1)}))$ and feedback $(p(h^{(k)}|h^{(k+1)}))$ paths. Denoting $x=h^{(0)}$ to simplify notation, the h update would thus consist in moves of the form

$$\begin{split} h^{(k)} \leftarrow & h^{(k)} + \delta \frac{\partial}{\partial h^{(k)}} \left(\log(p(h^{(k-1)}|h^{(k)})p(h^{(k)}|h^{(k+1)})) \right. \\ & + \left. \alpha \log(q(h^{(k)}|h^{(k-1)})q(h^{(k+1)}|h^{(k)})) \right), \end{split}$$

where α is as in Eq. 6. No back-propagation is needed for the above derivatives when $h^{(k)}$ is on the left hand side of the conditional probability bar. Sec. 6 deals with the right hand side case. For the left hand side case, e.g., $p(h^{(k)}|h^{(k+1)})$ a conditional Gaussian with mean μ and variance σ^2 , the gradient with respect to $h^{(k)}$ is simply $\frac{\mu-h^{(k)}}{\sigma^2}$. Note that there is an interesting interpretation of such a deep model: the layers above $h^{(k)}$ provide a complex implicitly defined prior for $p(h^{(k)})$.

5. Alternative Interpretations as Denoising Auto-Encoder

By inspection of Algorithm 1, one can observe that this algorithm trains p(x|h) and q(h|x) to form complementary pairs of an *autoencoder* (since the input of one is the target of the other and viceversa). Note that from that point of view any of the two can act as encoder and the other as decoder for it, depending on whether we start from h or from x. In the case of multiple latent layers, each pair of conditionals $q(h^{(k+1)}|h^{(k)})$ and $p(h^{(k)}|h^{(k+1)})$ forms a symmetric auto-encoder, i.e., either one can act as the encoder and the other as the corresponding decoder, since they are trained with the same $(h^{(k)}, h^{(k+1)})$ pairs (but with reversed roles of input and target).

In addition, if noise is injected, e.g., in the form of the quantization induced by a spike train, then the trained auto-encoders are actually denoising auto-encoders, which means that both the encoders and decoders are *contractive*: in the neighborhood of the observed (x,h) pairs, they map neighboring "corrupted" values to the "clean" (x,h) values.

5.1. Joint Denoising Auto-Encoder with Latent Variables

This suggests considering a special kind of "joint" denoising autoencoder which has the pair (x, h) as "visible" variable, an autoencoder that implicitly estimates an underlying p(x, h). The

transition operator³ for that joint visible-latent denoising autoencoder is the following in the case of a single hidden layer:

$$(\tilde{x}, \tilde{h}) \leftarrow \operatorname{corrupt}(x, h)$$

$$h \sim q(h|\tilde{x}) \qquad x \sim p(x|\tilde{h}), \tag{9}$$

where the corruption may correspond to the stochastic quantization induced by the neuron non-linearity and spiking process. In the case of a middle layer $h^{(k)}$ in a deeper model, the transition operator must account for the fact that $h^{(k)}$ can either be reconstructed from above or from below, yielding, with probability say $\frac{1}{2}$,

$$h^{(k)} \sim p(h^{(k)}|\tilde{h}^{(k+1)}),$$
 (10)

and with one minus that probability,

$$h^{(k)} \sim q(h^{(k)}|\tilde{h}^{(k-1)}).$$
 (11)

Since this interpretation provides a different model, it also provides a different way of *generating* samples. Especially for shallow, we have found that better samples could be obtained in this way, i.e., running the Markov chain with the above transition operator for a few steps.

There might be a geometric interpretation for the improved quality of the samples when they are obtained in this way, compared to the directed generative model that was defined earlier. Denote $q^*(x)$ the empirical distribution of the data, which defines a joint $q^*(h,x) = q^*(x)q^*(h|x)$. Consider the likely situation where p(x, h) is not well matched to $q^*(h, x)$ because for example the parametrization of p(h) is not powerful enough to capture the complex structure in the empirical distribution $q^*(h)$ obtained by mapping the training data through the encoder and inference $q^*(h|x)$. Typically, $q^*(x)$ would concentrate on a manifold and the encoder would not be able to completely unfold it, so that $q^*(h)$ would contain complicated structure with pockets or manifolds of high probability. If p(h) is a simple factorized model, then it will generate values of h that do not correspond well to those seen by the decoder p(x|h) when it was trained, and these out-of-manifold samples in h-space are likely to be mapped to out-of-manifold samples in x-space. One solution to this problem is to increase the capacity of p(h) (e.g., by adding more layers on top of h). Another is to make q(h|x) more powerful (which again can be achieved by increasing the depth of the model, but this time by inserting additional layers below h). Now, there is a cheap way of obtaining a very deep directed graphical model, by unfolding the Markov chain of an MCMC-based generative model for a fixed number of steps, i.e., considering each step of the Markov chain as an extra "layer" in a deep directed generative model, with shared parameters across these layers. As we have seen that there is such an interpretation via the joint denoising auto-encoder over both latent and visible, this idea can be immediately applied. We know that each step of the Markov chain operator moves its input distribution closer to the stationary distribution of the chain. So if we start from samples from a very broad (say factorized) prior p(h) and we iteratively encode/decode them (injecting noise appropriately as during training) by successively sampling from p(x|h) and then from q(h|x), the resulting h samples should end up looking more like those seen during training (i.e., from $q^*(h)$).

³See Theorem 1 from Bengio *et al.* (2013) for the generative interpretation of denoising auto-encoders: it basically states that one can sample from the model implicitly estimated by a denoising auto-encoder by simply alternating noise injection (corruption), encoding and decoding, these forming each step of a generative Markov chain.

5.2. Latent Variables as Corruption

There is another interpretation of the training procedure, also as a denoising auto-encoder, which has the advantage of producing a generative procedure that is the same as the inference procedure except for \boldsymbol{x} being unclamped.

We return again to the generative interpretation of the denoising criterion for auto-encoders, but this time we consider the non-parametric process $q^*(h|x)$ as a kind of corruption of x that yields the h used as input for reconstructing the observed x via p(x|h). Under that interpretation, a valid generative procedure consists at each step in first performing inference, i.e., sampling h from $q^*(h|x)$, and second sampling from p(x|h). Iterating these steps generates x's according to the Markov chain whose stationary distribution is an estimator of the data generating distribution that produced the training x's (Bengio $et\ al.$, 2013). This view does not care about how $q^*(h|x)$ is constructed, but it tells us that if p(x|h) is trained to maximize reconstruction probability, then we can sample in this way from the implicitly estimated model.

We have also found good results using this procedure (Algorithm 2 below), and from the point of view of biological plausibility, it would make more sense that "generating" should involve the same operations as "inference", except for the input being observed or not.

6. Targetprop instead of Backprop

In Algorithm 1 and the related stochastic variants Eq. 8 suggests that back-propagation (through one layer) is still needed when $h^{(k)}$ is on the right hand side of the conditional probability bar, e.g., to compute $\frac{\partial p(h^{(k-1)}|h^{(k)})}{\partial h^{(k)}}$. Such a gradient is also the basic building block in back-propagation for supervised learning: we need to back-prop through one layer, e.g. to make $h^{(k)}$ more "compatible" with $h^{(k-1)}$. This provides a kind *error signal*, which in the case of unsupervised learning comes from the sensors, and in the case of supervised learning, comes from the layer holding the observed "target".

Based on recent theoretical results on denoising auto-encoders, we propose the following estimator (up to a scaling constant) of the required gradient, which is related to previous work on "target propagation" (Bengio, 2014; Lee et~al., 2014) or targetprop for short. To make notation simpler, we focus below on the case of two layers h and x with "encoder" q(h|x) and "decoder" p(x|h), and we want to estimate $\frac{\partial \log p(x|h)}{\partial h}$. We start with the special case where p(x|h) is a Gaussian with mean g(h) and q(h|x) is Gaussian with mean f(x), i.e., f and g are the deterministic components of the encoder and decoder respectively. The proposed estimator is then

$$\widehat{\Delta h} = \frac{f(x) - f(g(h))}{\sigma_h^2},\tag{12}$$

where σ_h^2 is the variance of the noise injected in q(h|x).

Let us now justify this estimator. Theorem 2 by Alain and Bengio (2013) states that in a denoising auto-encoder with reconstruction function $r(x) = \operatorname{decode}(\operatorname{encode}(x))$, a well-trained auto-encoder estimates the log-score via the difference between its reconstruction and its input:

$$\frac{r(x) - x}{\sigma^2} \to \frac{\partial \log p(x)}{\partial x},$$

where σ^2 is the variance of injected noise, and p(x) is the implicitly estimated density. We are now going to consider two denoising auto-encoders and apply this theorem to them. First, we

note that the gradient $\frac{\partial \log p(x|h)}{\partial h}$ that we wish to estimate can be decomposed as follows:

$$\frac{\partial \log p(x|h)}{\partial h} = \frac{\partial \log p(x,h)}{\partial h} - \frac{\partial \log p(h)}{\partial h}.$$

Hence it is enough to estimate $\frac{\partial \log p(x,h)}{\partial h}$ as well as $\frac{\partial \log p(h)}{\partial h}$. The second one can be estimated by considering the auto-encoder which estimates p(h) implicitly and for which g is the encoder (with g(h) the "code" for h) and f is the decoder (with f(g(h))the "reconstruction" of h). Hence we have that $\frac{f(g(h))-h}{\sigma_h^2}$ is an estimator of $\frac{\partial \log p(h)}{\partial h}$.

The other gradient can be estimated by considering the joint denoising auto-encoder over (x, h) introduced in the previous section. The (noise-free) reconstruction function for that autoencoder is

$$r(x,h) = (g(h), f(x)).$$

Hence $\frac{f(x)-h}{\sigma_h^2}$ is an estimator of $\frac{\partial \log p(x,h)}{\partial h}$. Combining the two estimators, we get $\frac{(f(x)-h)}{\sigma_h^2} - \frac{(f(g(h))-h)}{\sigma_h^2} = \frac{f(x)-f(g(h))}{\sigma_h^2},$

$$\frac{(f(x)-h)}{\sigma_h^2} - \frac{(f(g(h))-h)}{\sigma_h^2} = \frac{f(x)-f(g(h))}{\sigma_h^2}$$

which corresponds to Eq. 12.

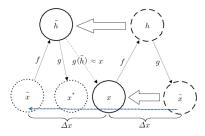


Figure 2. The optimal h for maximizing p(x|h) is \tilde{h} s.t. $g(\tilde{h}) =$ x. Since the encoder f and decoder g are approximate inverses of each other, their composition makes a small move Δx . Eq. 12 is obtained by assuming that by considering an \tilde{x} at $x - \Delta$ and applying $f \circ q$, one would approximately recover x, which should be true if the changes are small and the functions smooth (see Lee and Bengio (2014) for a detailed derivation).

Another way to obtain the same formula from a geometric perspective is illustrated in Figure 2. It was introduced in Lee and Bengio (2014) in the context of a backprop-free algorithm for training a denoising auto-encoder.

7. Related Work

An important inspiration for the proposed framework is the biological implementation of back-propagation proposed by Hinton (2007). In that talk, Hinton suggests that STDP corresponds to a gradient update step with the gradient on the voltage potential corresponding to its temporal derivative. To obtain the supervised back-propagation update in the proposed scenario would require symmetric weights and synchronization of the computations in terms of feedforward and feedback phases.

Our proposal introduces a novel machine learning interpretation that also matches the STDP behavior, based on a variational EM framework, allowing us to obtain a more biologically plausible mechanism for deep generative unsupervised learning, avoiding the need for symmetric weights, and introducing a novel method to obtain neural updates that approximately propagate gradients and move towards better overall configurations of neural activity (with difference target-prop). There is also an interesting connection with an earlier proposal for a more biologically plausible implementation of supervised back-propagation (Xie and Seung, 2003) which also relies on iterative inference (a deterministic relaxation in that case), but needs symmetric weights.

Another important inspiration is Predictive Sparse Decomposition (PSD) (Kavukcuoglu et al., 2008). PSD is a special case of Algorithm 1 when there is only one layer and the encoder q(h|x), decoder p(x|h), and prior p(h) have a specific form which makes p(x,h) a sparse coding model and q(h|x) a fast parametric approximation of the correct posterior. Our proposal extends PSD by providing a justification for the training criterion as a variational bound, by generalizing to multiple layers of latent variables, and by providing associated generative procedures.

The combination of a parametric approximate inference machine (the encoder) and a generative decoder (each with possibly several layers of latent variables) is an old theme that was started with the Wake-Sleep algorithm (Hinton et al., 1995) and which finds very interesting instantiations in the variational auto-encoder (Kingma and Welling, 2014; Kingma et al., 2014) and the reweighted wakesleep algorithm (Bornschein and Bengio, 2014). Two important differences with the approach proposed here is that here we avoid back-propagation thanks to an inference step that approximates the posterior. In this spirit, see the recent work introducing MCMC inference for the variational auto-encoder Salimans et al.

The proposal made here also owes a lot to the idea of target propagation introduced in Bengio (2014); Lee et al. (2014), to which it adds the idea that in order to find a target that is consistent with both the input and the final output target, it makes sense to perform iterative inference, reconciling the bottom-up and top-down pressures. Addressing the weight transport problem (the weight symmetry constraint) was also done for the supervised case using feedback alignment (Lillicrap et al., 2014): even if the feedback weights do not exactly match the feedforward weights, the latter learn to align to the former and "back-propagation" (with the wrong feedback weights) still works.

The targetprop formula avoiding back-propagation through one layer is actually the same as proposed by Lee and Bengio (2014) for backprop-free auto-encoders. What has been added here is a justification of this specific formula based on the denoising autoencoder theorem from Alain and Bengio (2013), and the empirical validation of its ability to climb the joint likelihood for variational inference.

Compared to previous work on auto-encoders, and in particular their generative interpretation (Bengio et al., 2013; 2014), this paper for the first time introduces latent variables without requiring back-propagation for training.

8. Experimental Validation

Figure 3 shows generated samples obtained after training on MNIST with Algorithm 2 (derived from the considerations of Sec.s 4, 5 and 6). The network has two hidden layers, h_1 with 1000 softplus units and h_2 with 100 sigmoid units (which can be considered biologically plausible (Glorot et al., 2011)). We trained for 20 epochs, with minibatches of size 100 to speedup computation using GPUs. Results can be reproduced from code at http://goo.gl/hoQqR5. Using the Parzen density estimator previously used for that data, we obtain an estimated log-likelihood LL=236 (using a standard deviation of 0.2

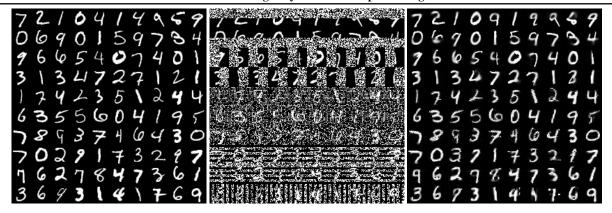


Figure 5. Examples of filling-in (in-painting) missing (initially corrupted) parts of an image. Left: original MNIST test examples. Middle: initial state of the inference, with half of the pixels randomly sampled (with a different corruption pattern in each row of the figure). Right: reconstructions using a variant of the INFERENCE procedure of Algorithm 2 for the case when some inputs are clamped.



Figure 3. MNIST samples generated by GENERATE from Algorithm 2 after training with TRAIN.

for the Parzen density estimator, chosen with the validation set), which is about the same or better as was obtained for contractive auto-encoders (Rifai et al., 2011) (LL=121), deeper generative stochastic networks (Bengio et al., 2014) (LL=214) and generative adversarial networks (Goodfellow et al., 2014) (LL=225). In accordance with Algorithm 2, the variances of the conditional densities are 1, and the top-level prior is ignored during most of training (as if it was a very broad, uniform prior) and only set to the Gaussian by the end of training, before generation (by setting the parameters of $p(h_2)$ to the empirical mean and variance of the projected training examples at the top level). Figure 4 shows that the targetprop updates (instead of the gradient updates) allow the inference process to indeed smoothly increase the joint likelihood. Note that if we sample using the directed graphical model p(x|h)p(h), the samples are not as good and LL=126, suggesting as discussed in Sec. 5 that additional inference and encode/decode iterations move h towards values that are closer to $q^*(h)$ (the empirical distribution of inferred states from training examples). The experiment illustrated in Figure 5 shows that the proposed inference mechanism can be used to fill-in missing values with a trained model. The model is the same that was trained using Algorithm 2 (with samples shown in Figure 3). 20 iterations steps of encode/decode as described below were performed, with a call

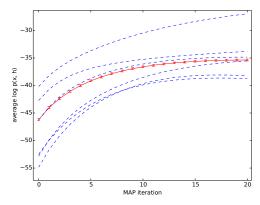


Figure 4. Increase of $\log p(x,h)$ over 20 iterations of the INFER-ENCE algorithm 2, showing that the targetprop updates increase the joint likelihood. The solid red line shows the average and the standard error over the full testset containing 10,000 digits. Dashed lines show $\log p(x,h)$ for individual datapoints.

to INFERENCE (to maximize p(x,h) over h) for each step, with a slight modification. Instead of using $f_1(x)-f_1(g_1(h))$ to account for the pressure of x upon h (towards maximizing p(x|h)), we used $f_1(x^v,g^m(h))-f(g(h))$, where x^v is the part of x that is visible (clamped) while $g^m(h)$ is the part of the output of g(h) that concerns the missing (corrupted) inputs. This formula was derived from the same consideration as for Eq. 12, but where the quantity of interest is $\frac{\partial \log p(x^v|h)}{\partial h}$ rather than $\frac{\partial \log p(x|h)}{\partial h}$, and we consider that the reconstruction of h, given x^v , fills-in the missing inputs (x^m) from $g^m(h)$.

9. Future Work and Conclusion

We consider this paper as an exploratory step towards explaining a central aspect of the brain's learning algorithm: *credit assignment through many layers*. Of the non-plausible elements of back-propagation described in the introduction, the proposed approach addresses all except the 5th. As argued by Bengio (2014); Lee *et al.* (2014), departing from back-propagation could be useful not just for biological plausibility but from a machine learning point of view as well: by working on the "targets" for the intermediate layers, we may avoid the kind of reliance on smoothness and derivatives that characterizes back-propagation, as these tech-

niques can in principle work even with highly non-linear transformations for which gradients are often near 0, e.g., with stochastic binary units (Lee *et al.*, 2014). Besides the connection between STDP and variational EM, an important contribution of this paper is to show that the "targetprop" update which estimates the gradient through one layer can be used for inference, yielding systematic improvements in the joint likelihood and allowing to learn a good generative model. Another interesting contribution is that the variational EM updates, with noise added, can also be interpreted as training a denoising auto-encoder over both visible and latent variables, and that iterating from the associated Markov chain yields better samples than those obtained from the directed graphical model estimated by variational EM.

Many directions need to be investigated to follow-up on the work reported here. An important element of neural circuitry is the strong presence of lateral connections between nearby neurons in the same area. In the proposed framework, an obvious place for such lateral connections is to implement the prior on the joint distribution between nearby neurons, something we have not explored in our experiments. For example, Garrigues and Olshausen (2008) have discussed neural implementations of the inference involved in sparse coding based on the lateral connections.

Although we have found that "injecting noise" helped training a better model, more theoretical work needs to be done to explore this replacement of a MAP-based inference by an MCMC-like inference, which should help determine how and how much of this noise should be injected.

Whereas this paper focused on unsupervised learning, these ideas could be applied to supervised learning and reinforcement learning as well. For reinforcement learning, an important role of the proposed algorithms is to learn to predict rewards, although a more challenging question is how the MCMC part could be used to simulate future events. For both supervised learning and reinforcement learning, we would probably want to add a mechanism that would give more weight to minimizing prediction (or reconstruction) error for some of the observed signals (e.g. y is more important to predict than x).

Finally, a lot needs to be done to connect in more detail the proposals made here with biology, including neural implementation using spikes with Poisson rates as the source of signal quantization and randomness, taking into account the constraints on the sign of the weights depending on whether the pre-synaptic neuron is inhibitory or excitatory, etc. In addition, although the operations proposed here are backprop-free, they may still require some kinds of synchronizations (or control mechanism) and specific connectivity to be implemented in brains.

Acknowledgments

The authors would like to thank Jyri Kivinen, Tim Lillicrap and Saizheng Zhang for feedback and discussions, as well as NSERC, CIFAR, Samsung and Canada Research Chairs for funding, and Compute Canada for computing resources.

References

Alain, G. and Bengio, Y. (2013). What regularized auto-encoders learn from the data generating distribution. In *ICLR* '2013. also arXiv report 1211.4246.

- Bengio, Y. (2014). How auto-encoders could provide credit assignment in deep networks via target propagation. Technical report, arXiv:1407.7906.
- Bengio, Y., Yao, L., Alain, G., and Vincent, P. (2013). Generalized denoising auto-encoders as generative models. In *NIPS'2013*.
- Bengio, Y., Thibodeau-Laufer, E., Alain, G., and Yosinski, J. (2014). Deep generative stochastic networks trainable by backprop. In *ICML*'2014.
- Bi, G. Q. and Poo, M. M. (1998). Synaptic modifications in cultured hippocampal neurons: dependence on spike timing, synaptic strength, and postsynaptic cell type. *J Neurosci*, 18, 10464–10472.
- Bornschein, J. and Bengio, Y. (2014). Reweighted wake-sleep. Technical report, arXiv:1406.2751.
- de Freitas, N., Højen-Sørensen, P., Jordan, M. I., and Russell, S. (2001). Variational MCMC. In *UAI'2001*, pages 120–127.
- Douglas, R. J., Martin, K. A., and Whitteridge, D. (1989). A canonical microcircuit for neocortex. *Neural Computation*, 1(480–488).
- Felleman, D. J. and Essen, D. C. V. (1991). Distributed hierarchical processing in the primate cerebral cortex. *Cerebral Cortex*, **1**, 1–47.
- Garrigues, P. and Olshausen, B. (2008). Learning horizontal connections in a sparse coding model of natural images. In J. Platt, D. Koller, Y. Singer, and S. Roweis, editors, *Advances in Neural Information Processing Systems 20 (NIPS'07)*, pages 505–512. MIT Press, Cambridge, MA.
- Gerstner, W., Kempter, R., van Hemmen, J., and Wagner, H. (1996). A neuronal learning rule for sub-millisecond temporal coding. *Nature*, **386**, 76–78.
- Glorot, X., Bordes, A., and Bengio, Y. (2011). Deep sparse rectifier neural networks. In *AISTATS* '2011.
- Goodfellow, I., Le, Q., Saxe, A., and Ng, A. (2009). Measuring invariances in deep networks. In *NIPS* '2009, pages 646–654.
- Goodfellow, I. J., Pouget-Abadie, J., Mirza, M., Xu, B., Warde-Farley, D., Ozair, S., Courville, A., and Bengio, Y. (2014). Generative adversarial networks. In NIPS'2014.
- Hinton, G. E. (2007). How to do backpropagation in a brain. Invited talk at the NIPS'2007 Deep Learning Workshop.
- Hinton, G. E. and Sejnowski, T. J. (1986). Learning and relearning in Boltzmann machines. In D. E. Rumelhart and J. L. McClelland, editors, *Parallel Distributed Processing: Explorations in the Microstructure of Cognition. Volume 1: Foundations*, pages 282–317. MIT Press, Cambridge, MA.
- Hinton, G. E., Dayan, P., Frey, B. J., and Neal, R. M. (1995). The wake-sleep algorithm for unsupervised neural networks. *Science*, **268**, 1558–1161.
- Hinton, G. E., Osindero, S., and Teh, Y. (2006). A fast learning algorithm for deep belief nets. *Neural Computation*, **18**, 1527–1554.

- Kavukcuoglu, K., Ranzato, M., and LeCun, Y. (2008). Fast inference in sparse coding algorithms with applications to object recognition. Technical report, Computational and Biological Learning Lab, Courant Institute, NYU. Tech Report CBLL-TR-2008-12-01.
- Kingma, D., Rezende, D., Mohamed, S., and Welling, M. (2014).Semi-supervised learning with deep generative models. In NIPS'2014.
- Kingma, D. P. and Welling, M. (2014). Auto-encoding variational bayes. In *Proceedings of the International Conference on Learning Representations (ICLR)*.
- Lee, D.-H. and Bengio, Y. (2014). Backprop-free auto-encoders. NIPS'2014 Deep Learning workshop.
- Lee, D.-H., Zhang, S., Biard, A., and Bengio, Y. (2014). Target propagation. arXiv 1412.7525.
- Legenstein, R. and Maass, W. (2014). Ensembles of spiking neurons with noise support optimal probabilistic inference in a dynamically changing environment. *PLOS Computational Biology*, **10**(10).
- Lillicrap, T. P., Cownden, D., Tweed, D. B., and Akerman, C. J. (2014). Random feedback weights support learning in deep neural networks. arXiv:1411.0247.
- Markram, H. and Sakmann, B. (1995). Action potentials propagating back into dendrites triggers changes in efficacy. *Soc. Neurosci. Abs*, **21**.
- Neal, R. and Hinton, G. (1999). A view of the EM algorithm that justifies incremental, sparse, and other variants. In M. I. Jordan, editor, *Learning in Graphical Models*. MIT Press, Cambridge, MA.
- Pecevski, D., Buesing, L., and Maass, W. (2011). Probabilistic inference in general graphical models through sampling in stochastic networks of spiking neurons. *PLOS Computational Biology*.
- Ranzato, M., Poultney, C., Chopra, S., and LeCun, Y. (2007). Efficient learning of sparse representations with an energy-based model. In NIPS'2006.
- Rezende, D. J. and Gerstner, W. (2014). Stochastic variational learning in recurrent spiking networks. Frontiers in Computational Neuroscience, 8(38).
- Rifai, S., Vincent, P., Muller, X., Glorot, X., and Bengio, Y. (2011). Contractive auto-encoders: Explicit invariance during feature extraction. In *ICML'2011*.
- Salakhutdinov, R. and Hinton, G. (2009). Deep Boltzmann machines. In AISTATS'2009.
- Salimans, T., Kingma, D. P., and Welling, M. (2014). Markov Chain Monte Carlo and Variational Inference: Bridging the Gap. arXiv:1410.6460.
- Sjöström, J. and Gerstner, W. (2010). Spike timing dependent plasticity. *Scholarpedia*, **5**(2).
- Smolensky, P. (1986). Information processing in dynamical systems: Foundations of harmony theory. In D. E. Rumelhart and J. L. McClelland, editors, *Parallel Distributed Processing*, volume 1, chapter 6, pages 194–281. MIT Press, Cambridge.

- Vincent, P., Larochelle, H., Bengio, Y., and Manzagol, P.-A. (2008). Extracting and composing robust features with denoising autoencoders. In *ICML* 2008.
- Welling, M. and Teh, Y.-W. (2011). Bayesian learning via stochastic gradient Langevin dynamics. In *ICML'2011*, pages 681–688.
- Williams, R. J. (1992). Simple statistical gradient-following algorithms connectionist reinforcement learning. *Machine Learning*, 8, 229–256.
- Xie, X. and Seung, H. S. (2000). Spike-based learning rules and stabilization of persistent neural activity. In S. Solla, T. Leen, and K. Müller, editors, *Advances in Neural Information Pro*cessing Systems 12, pages 199–208. MIT Press.
- Xie, X. and Seung, H. S. (2003). Equivalence of backpropagation and contrastive Hebbian learning in a layered network. *Neural Computation*.

Algorithm 2 Inference, training and generation procedures used in Experiment 1. The algorithm can naturally be extended to more layers. $f_i()$ is the feedforward map from layer i-1 to layer i and $g_i()$ is the feedback map from layer i to layer i-1, with $x=h_0$ being layer 0.

```
Define INFERENCE(x, N=15, \delta=0.1, \alpha=0.001):
Feedforward pass: h_1 \leftarrow f_1(x), h_2 \leftarrow f_2(h_1)
\mathbf{for}\ t = 1\ \mathbf{to}\ N\ \mathbf{do}
   h_2 \leftarrow h_2 + \delta(f_2(h_1) - f_2(g_2(h_2)))
   h_1 \leftarrow h_1 + \delta(f_1(x) - f_1(g_1(h_1))) + \alpha(g_2(h_2) - h_1)
end for
Return h_1, h_2
Define TRAIN()
for x in training set do
   h_1, h_2 \leftarrow \text{INFERENCE}(x) (i.e. E-part of EM)
   update each layer using local targets (M-part of EM) \Theta \leftarrow \Theta + \epsilon \frac{\partial}{\partial \Theta} \left( g_l(\tilde{h}_l) - h_{l-1})^2 \right. \\ \Theta \leftarrow \Theta + \epsilon \frac{\partial}{\partial \Theta} \left( f_l(\tilde{h}_{l-1}) - h_l \right)^2
    where \tilde{h}_l is a Gaussian-corrupted version of h_l. For
    the top sigmoid layer \tilde{h}_2 we average 3 samples from a
    Bernoulli distribution with p(h_2 = 1) = h_2 to obtain
   a spike-like corruption.
end for
```

Compute the mean and variance of h_2 using the training set. Multiply the variances by 4. Define $p(h_2)$ as sampling from this Gaussian.

```
Define GENERATE(): Sample h_2 from p(h_2) Assign h_1 \leftarrow g_2(h_2) and x \leftarrow g_1(h_1) for t=1 to 3 do h_1, h_2 \leftarrow \text{INFERENCE}(x, \alpha=0.3) x \leftarrow g_1(h_1) end for Return x
```