UNDERSTANDING SELF-SUPERVISED LEARNING WITH DUAL DEEP NETWORKS

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ABSTRACT

We propose a novel theoretical framework to understand self-supervised learning methods that employ dual pairs of deep ReLU networks (e.g., SimCLR, BYOL). First, we prove that in each SGD update of SimCLR, the weights at each layer are updated by a covariance operator that specifically amplifies initial random selectivities that vary across data samples but survive averages over data augmentations, which we show leads to the emergence of hierarchical features, if the input data are generated from a hierarchical latent tree model. With the same framework, we also show analytically that in BYOL, the usage of BatchNorm and a predictor creates an implicit contrastive term, acting as an approximate covariance operator. The term is formed by the inter-play between the zero-mean operation of BatchNorm and the extra predictor in the online network. Extensive ablation studies justify our theoretical findings.

1 INTRODUCTION

While self-supervised learning (SSL) has achieved great empirical success across multiple domains, including computer vision (He et al., 2020; Goyal et al., 2019; Chen et al., 2020a; Grill et al., 2020; Misra and Maaten, 2020; Caron et al., 2020), natural language processing (Devlin et al., 2018), and speech recognition (Wu et al., 2020; Baevski and Mohamed, 2020; Baevski et al., 2019), its theoretical understanding remains elusive, especially when multi-layer nonlinear deep networks are involved. Unlike supervised learning (SL) that deals with labeled data, SSL learns meaningful structures from randomly initialized networks without human-provided labels.

In this paper, we propose a systematic theoretical analysis of SSL with deep ReLU networks. Our analysis imposes no parametric assumptions on the input data distribution and is applicable to state-of-the-art SSL methods that typically involve two parallel (or dual) deep ReLU networks during training (e.g., SimCLR (Chen et al., 2020a), BYOL (Grill et al., 2020), etc). We do so by developing an analogy between SSL and a theoretical framework for analyzing supervised learning, namely the student-teacher setting (Tian, 2020; Allen-Zhu and Li, 2020; Lampinen and Ganguli, 2018; Saad and Solla, 1996), which also employs a pair of dual networks. Our results indicate that SimCLR weight updates at every layer are amplified by a fundamental positive semi definite (PSD) covariance operator that only captures feature variability across data points that survive averages over data augmentation procedures designed in practice to scramble semantically unimportant features (e.g. random image crops, blurring or color distortions (Falcon and Cho, 2020; Kolesnikov et al., 2019; Misra and Maaten, 2020; Purushwalkam and Gupta, 2020)). This covariance operator provides a principled framework to study how SimCLR amplifies initial random selectivity to obtain distinctive features that vary across samples after surviving averages over data-augmentations.

Based on the covariance operator, we further show that (1) in a two-layer setting, a top-level covariance operator helps accelerate the learning of low-level features, and (2) when the data are generated by a hierarchical latent tree model, training deep ReLU networks leads to an emergence of the latent variables in its intermediate layers. We also tentatively explain why BYOL works without negative pairs, by showing analytically that an interplay between the zero-mean operation in BatchNorm and the extra predictor in the online network creates an implicit contrastive term, which is consistent with empirical observations in the recent blog (Fetterman and Albrecht, 2020). Note that our “sufficient

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Figure 1: (a) Overview of the two SSL algorithms we study in this paper: SimCLR (\(W_1 = W_2 = W\); no predictor, NCE Loss) and BYOL (\(W_2\) has an extra predictor, \(W_2\) is a moving average). (b) Detailed Notation.

reasoning” does not rule out other normalization techniques that also make BYOL work without negative pairs, as pointed by a recent update of BYOL (Richemond et al., 2020). Finally, we also discover that reinitializing the predictor every a few epochs doesn’t hurt the performance measured by linear evaluation protocol, questioning the hypothesis of “optimal predictor” in a recent version (v3) of BYOL (Grill et al., 2020).

To the best of our knowledge, we are the first to provide a systematic theoretical analysis of modern SSL methods with deep ReLU networks that also incorporates data augmentation, and show how the gradient-based SSL work inside the neural network.

2 RELATED WORKS

In addition to SimCLR and BYOL, many concurrent SSL frameworks exist to learn good representations for computer vision tasks. MoCo (He et al., 2020; Chen et al., 2020b) keeps a large bank of past representations in a queue as the slow-progressing target to train from. DeepCluster (Caron et al., 2018) and SwAV (Caron et al., 2020) learn the representations by iteratively or implicitly clustering on the current representations and improving representations using the cluster label. (Alwassel et al., 2019) applies similar ideas to multi-modality tasks. Contrastive Predictive Coding (Oord et al., 2018) learns the representation by predicting the future of a sequence in the latent space with autoregressive models and InfoNCE loss. Contrastive MultiView Coding (Tian et al., 2019) uses multiple sensory channels (called “view”) of the same scene as the positive pairs and independently sampled views as the negative pairs to train the model. Recently, (Li et al., 2020) moves beyond instance-wise pairs and proposes to use prototypes to construct training pairs that are more semantically meaningful.

In contrast, the literature on the (theoretical) analysis of SSL is sparse. (Wang and Isola, 2020) shows directly optimizing the alignment/uniformity of the positive/negative pairs leads to comparable performance against contrastive loss. (Arora et al., 2019a) proposes an interesting analysis of how contrastive learning aids downstream classification tasks, given assumptions about data generation. (Lee et al., 2020) analyzes how learning pretext tasks could reduce the sample complexity of the downstream task and (Tosh et al., 2020) analyzes contrastive loss with multi-view data in the semi-supervised setting, with different generative models. However, they either work on linear model or treat deep models as a black-box function approximator with sufficient capacity. In comparison, we incorporate self-supervision, deep models, contrastive loss, data augmentation and generative models together into the same theoretical framework, and makes an attempt to understand how the intermediate features are emerged in modern SSL architectures with deep models that achieve SoTA.

3 OVERALL FRAMEWORK

Notation. Consider an \(L\)-layer ReLU network obeying \(f_l = \psi(f_{l-1})\) and \(\hat{f}_l = W_l f_{l-1}\) for \(l = 1, \ldots, L\). Here \(f_0 = x\) being the input, and \(f_L = f\) the output (no ReLU at the top layer). \(W_l \in \mathbb{R}^{n_l \times n_{l-1}}\) are the weight matrices, and \(\psi(u) := \max(u, 0)\) is the element-wise ReLU nonlinearity. We let \(\mathcal{W} := \{W_l\}_{l=1}^L\) be all network weights. We also denote the gradient of any loss function with respect to \(f_l\) by \(g_l \in \mathbb{R}^{n_l}\), and the derivative of the output \(f_L\) with respect to an earlier pre-activation \(f_l\) by the Jacobian matrix \(J_l(x; \mathcal{W}) \in \mathbb{R}^{n_L \times n_l}\), as both play key roles in backpropagation (Fig. 1(b)).

An analogy between self-supervised and supervised learning: the dual network scenario. Many recent successful approaches to self-supervised learning (SSL), including SimCLR (Chen
et al., 2020a), BYOL (Grill et al., 2020) and MoCo (He et al., 2020), employ a dual “Siamese-like” pair (Koch et al., 2015) of such networks (Fig. 1(b)). Each network has its own set of weights $W_1$ and $W_2$, receives respective inputs $x_1$ and $x_2$ and generates outputs $f_{1,L}(x_1;W_1)$ and $f_{2,L}(x_2;W_2)$. The pair of inputs $\{x_1, x_2\}$ can be either positive or negative, depending on how they are sampled. For a positive pair, a single data point $x$ is drawn from the data distribution $p(\cdot)$, and then two augmented views $x_1$ and $x_2$ are drawn from a conditional augmentation distribution $p_{\text{aug}}(\cdot|x)$. Possible image augmentations include random crops, blurs or color distortions, that ideally preserve semantic content useful for downstream tasks. In contrast, for a negative pair, two different data points $x,x' \sim p(\cdot)$ are sampled, and then each are augmented independently to generate $x_1 \sim p_{\text{aug}}(\cdot|x)$ and $x_2 \sim p_{\text{aug}}(\cdot|x')$. For SimCLR, the dual networks have tied weights with $W_1 = W_2$, and a loss function is chosen to encourage the representation of positive (negative) pairs to become similar (dissimilar). In BYOL, only positive pairs are used, and the first network $W_1$, called the online network, is trained to match the output of the second network $W_2$ (the target), using an additional layer named predictor. The target network ideally provides training targets that can improve the online network’s representation and does not contribute gradient. The improved online network is gradually incorporated into the target network, yielding a bootstrapping procedure.

Our fundamental goal is to analyze the mechanisms governing how SSL methods like SimCLR and BYOL lead to the emergence of meaningful intermediate features, starting from random initializations, and how these features depend on the data distribution $p(x)$ and augmentation procedure $p_{\text{aug}}(\cdot|x)$. Interestingly, the analysis of supervised learning (SL) often employs a similar dual network scenario, called teacher-student setting (Tian, 2020; Allen-Zhu and Li, 2020; Lampinen and Ganguli, 2018; Saad and Solla, 1996), where $W_2$ are the ground truth weights of a fixed teacher network, which generates outputs in response to random inputs. These input-output pairs constitute training data for the first network, which is a student network. Only the student network’s weights $W_1$ are trained to match the target outputs provided by the teacher. This yields an interesting mathematical parallel between SL, in which the teacher is fixed and only the student evolves, and SSL, in which both the teacher and student evolve with potentially different dynamics. This mathematical parallel opens the door to using techniques from SL (e.g., (Tian, 2020)) to analyze SSL.

**Gradient of $\ell_2$ loss for dual deep ReLU networks.** As seen above, the (dis)similarity of representations between a pair of dual networks plays a key role in both SSL and SL. We thus consider minimizing a simple measure of dissimilarity, the squared $\ell_2$ distance $r := \frac{1}{2}\|f_{1,L} - f_{2,L}\|^2$ between the final outputs $f_{1,L}$ and $f_{2,L}$ of two multi-layer ReLU networks with weights $W_1$ and $W_2$ and inputs $x_1$ and $x_2$. Without loss of generality, we only analyze the gradient w.r.t $W_1$. For each layer $l$, we first define connection $K_l(x)$, a quantity that connects the bottom-up feature vector $f_{l-1}$ with the top-down Jacobian $J_l$, which both contribute to the gradient at weight layer $l$.

**Definition 1** (The connection $K_l(x)$). The connection $K_l(x;W) := f_{l-1}(x;W) \otimes J_l^T(x;W) \in \mathbb{R}^{n_{l-1} \times n_l}$. Here $\otimes$ is the Kronecker product.

**Theorem 1** (Squared $\ell_2$ Gradient for dual deep ReLU networks). The gradient $g_{W_1}$ of $r$ w.r.t. $W_1 \in \mathbb{R}^{n_1 \times n_1}$ for a single input pair $\{x_1, x_2\}$ is (here $K_{1,l} := K_1(x_1;W_1)$ and $K_{2,l} := K_1(x_2;W_2)$):

$$g_{W_1} = \text{vec} \left( \frac{\partial r}{\partial W_{1,l}} \right) = K_{1,l} \left[ K_{1,l}^T \text{vec}(W_{1,l}) - K_{2,l}^T \text{vec}(W_{2,l}) \right].$$

We used vectorized notation for the gradient $g_{W_1}$ and weights $W_1$ to emphasize certain theoretical properties of SSL learning below. The equivalent matrix form is $\frac{\partial r}{\partial W_{1,l}} = J_l^T [J_{1,l} W_{1,l} J_{1,l-1}^T - J_{2,l} W_{2,l} J_{2,l-1}^T] f_{l-1,l}^T$. See Appendix for proofs of all theorems in main text.

**Remarks.** Note that Theorem 1 can also be applied to networks with other non-linearity than ReLU, as long as they satisfy a more general property called reversibility (See Def. 2 in Appendix). This includes linear, LeakyReLU, and monomial activations, such as the quadratic activation function $\psi(x) = x^2$ used in many previous works (Du and Lee, 2018; Soltanolkotabi et al., 2018; Allen-Zhu and Li, 2020). In this paper, we focus on ReLU and leave other non-linearity for future works.

4 ANALYSIS OF SIMCLR

As discussed above, SimCLR (Chen et al., 2020a) employs both positive and negative input pairs, and a symmetric network structure with $W_1 = W_2 = \mathcal{W}$. Let $\{x_1, x_2\}$ be a positive input pair, and let $\{x_k\}$ for $k = 1, \ldots, H$ be $H$ negative pairs. These input pairs induce corresponding squared $\ell_2$ distances in output space, $r_+ := \frac{1}{2} \|f_{1,L} - f_{+,L}\|^2$, and $r_{k-} := \frac{1}{2} \|f_{1,L} - f_{k-,L}\|^2$. 

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The InfoNCE loss (Oord et al., 2018) with temperature $\tau$ then minimizes (maximizes) the positive (negative) pair distances:

$$L(r_+, r_{-1}, r_{-2}, \ldots, r_{-H}) := -\log \frac{e^{-r_+/\tau}}{e^{-r_+/\tau} + \sum_{k=1}^{H} e^{-r_{-k}/\tau}}$$  \hspace{1cm} (2)

When $\|u\|_2 = \|v\|_2 = 1$, we have $-\frac{1}{2}\|u - v\|_2^2 = \text{sim}(u, v) - 1$ where $\text{sim}(u, v) = \frac{u^T v}{\|u\|_2\|v\|_2}$.

and Eqn. 2 reduces to what the original SimCLR uses (the term $e^{-1/\tau}$ cancels out).

Note that in our analysis, for now we haven’t incorporated $\ell_2$ normalization at the topmost layer. The BYOL paper (Appendix F.6 in Grill et al. (2020) v3) shows that even without $\ell_2$ normalization, the algorithm still works despite numerical instabilities. Without normalization, the goal of analysis is to show that useful weight components grow exponential after gradient update (and we assume that the extra normalization will stabilize it).

One property of InfoNCE is important for our analysis:

**Theorem 2 (Property of Contrastive Loss).** For InfoNCE loss (Eqn. 2), we have $\frac{\partial L}{\partial r_k} > 0$ and $\frac{\partial L}{\partial r_k} < 0$ for $k = 1, \ldots, H$ and $\frac{\partial L}{\partial r_k} + \sum_{k=1}^{H} \frac{\partial L}{\partial r_{k-1}} = 0$.

Now, under an approximation in which we neglect variations in $\frac{\partial L}{\partial r_k}$ across different $k$ in a large minibatch, Theorem 1 and 2 imply the SimCLR gradient is governed by positive semi-definite (PSD) covariance operator at any layer (some simple contrastive loss (e.g., $r_+ - r_-$) satisfies it trivially):

**Theorem 3 (Covariance Operator from Contrastive Loss).** Assume that for any $\{x_1, x_+, x_{k-}\}$, $\frac{\partial L}{\partial r_k} = -\beta/H$ is a constant, then in a large batch limit, the update for $W_l$ is

$$W_l(t+1) = W_l(t) + \alpha \Delta W_l(t), \quad \text{where vec}(\Delta W_l(t)) = \beta V_x [K_l(x)] \text{vec}(W_l(t))$$  \hspace{1cm} (3)

where $\alpha$ is the learning rate and $K_l(x; W) := E_{x' \sim p_{\text{aug}}(\cdot|x)} [K_l(x'; W)]$ is the expected connection under the augmentation distribution, conditional on the datapoint $x$.

Above, we use the notation $\text{Cov}[X, Y] := E[XY^\top] - E[X]E[Y]^\top$ and $\forall [X] := \text{Cov}[X, X]$. The covariance operator $V_x [K_l(x)] \in \mathbb{R}^{n_l \times n_l \times n_l \times n_l}$ is a time-varying PSD matrix over the entire training procedure. Therefore, all its eigenvalues are non-negative and at any time $t$, $W_l$ is most amplified along its largest eigenmodes. Intuitively, this covariance operator ignores different views of the same sample $x$ by averaging over the augmentation distribution to compute $K_l(x)$, and then computes the expected covariance of this augmentation averaged connection with respect to the data distribution $p(x)$. Thus, at all layers, any variability in the connection across different data points, that survives augmentation averages, leads to weight amplification. This amplification of weights by the PSD data covariance of an augmentation averaged connection constitutes a fundamental description of SimCLR learning dynamics for arbitrary data and augmentation distributions.

4.1 DISCUSSIONS

**Role played by Theorem 2.** In the proof of Theorem 3, the condition $\frac{\partial L}{\partial r_k} + \sum_{k=1}^{H} \frac{\partial L}{\partial r_{k-1}} = 0$ from Theorem 2 is used to cancel out the term $K_l(x_1)K_l^\top(x_1)$, which becomes $E_x [K_l(x)K_l^\top(x)]$ in the large batch limit. Note that this expectation is with respect to all variations of $x$, including sampling from the dataset $p(\cdot)$ and data augmentation using $p_{\text{aug}}(\cdot|x)$. Therefore, this extra term neither encouraging separation of representations of different samples, nor enforcing invariance of representations within data augmentation, and is not likely to be useful. By using InfoNCE loss, SimCLR implicitly eliminates this unnecessary term. We haven’t seen previous works that discover this property and use it to analyze.

**The assumption of constant $\frac{\partial L}{\partial r_k}$**. This condition is mainly used technically to simplify the proof to yield a clean conclusion. We will leave its extension to non-constant $\frac{\partial L}{\partial r_k}$ in our future work. Note that if the contrastive loss is simply $r_+ - r_-$ then this condition trivially holds. For InfoNCE, $\frac{\partial L}{\partial r_k}$ and $\frac{\partial L}{\partial r_{k-1}}$ vary over tuples $(x_1, x_+, x_{k-})$, which creates complicated interactions among these samples. A high-level conceptual modeling is that InfoNCE puts emphasize on those “most violated” tuples that have similar representations for negative pairs and dis-similar representations for positive pairs. As a result, the actual covariance operator for SimCLR without the constant assumption is somehow a weighted version of $V_x [K_l(x)]$. 


Difference from Neural Tangent Kernels (NTK). Note that our covariance operator is a completely different mathematical object than Neural Tangent Kernel (NTK) (Jacot et al., 2018; Arora et al., 2019b). NTK is defined in the sample space and is full-rank if samples are distinct. For super wide networks (compared to the sample size), NTK seldom changes during training and leads to a convex optimization landscape. On the other hand, the covariance operator $\mathbb{V}_x[K_l(x)]$ is defined per layer on any data distribution and networks with finite width, and does not grow in size with larger sample size. $\mathbb{V}_x[K_l(x)]$ may change over the entire training procedure but remains PSD, leading to many interesting behaviors beyond learning coefficients on kernel feature space (e.g., as shown in Sec. 5.2, under the covariance operator, dynamics of high-layer weights leads to faster low-layer training). Furthermore, while NTK has nothing to do with data augmentation, $\mathbb{V}_x[K_l(x)]$ is tied to data augmentation and SSL architectures.

5 HOW THE COVARIANCE OPERATOR DRIVES THE EMERGENCE OF FEATURES

To concretely illustrate how the fundamental covariance operator derived in Theorem 3 drives feature emergence in SimCLR, we setup the following paradigm for analysis. We assume the input $x = x(z_0, z')$ is generated by two groups of latent variables, class/sample-specific latent $z_0$ and nuisance latent $z'$. After data augmentation, $z_0$ remains the same while $z'$ changes (Fig. 2(a)). In this case, the covariance operator is simply $\mathbb{V}_{z_0}[K_l(z_0)]$ since all nuisance latents $z'$ are integrated out in $K_l(z_0)$.

In this setting, we first show that a linear neuron performs PCA within an augmentation preserved subspace. We then consider how nonlinear neurons with local receptive fields (RFs) can learn to detect simple objects. Finally, we extend our analysis to deep ReLU networks exposed to data generated by a hierarchical latent tree model (HLTM), proving that, with sufficient over-parameterization, there exist lucky nodes at initialization whose activation is correlated with latent variables underlying the data, and that SimCLR amplifies these initial lucky patterns during learning.

### 5.1 Self-supervised Learning and the Single Neuron: Illustrative Examples

**A single linear neuron performs PCA in a preserved subspace.** For a single linear neuron ($L = 1, n_L = 1$), the connection in definition 1 is simply $K_1(x) = x$. Now imagine the input space $x$ can be decomposed into the direct sum of a semantically relevant subspace, and its orthogonal complement, which corresponds to a subspace of nuisance features. Furthermore, suppose the augmentation distribution $p_{aug}(|x|)$ is obtained by multiplying $x$ by a random Gaussian matrix that acts only in the nuisance subspace, thereby identically preserving the semantic subspace. Then the augmented averaged connection $\bar{K}_1(x) = Q^s x$ where $Q^s$ is a projection operator onto the semantic subspace. In essence, only the projection of data onto the semantic subspace survives augmentation averaging, as the nuisance subspace is scrambled. Then the covariance operator in Theorem 3 is $\mathbb{V}_x[K_1(x)] = Q^s \mathbb{V}_x[x] Q^{sT}$. Thus the covariance of the data distribution, projected onto the semantic subspace, governs the growth of the weight vector $W_1$, demonstrating SimCLR on a single linear neuron performs PCA within a semantic subspace preserved by data augmentation.

**A single linear neuron cannot detect localized objects.** We now consider a generative model in which data vectors can be thought of as images of objects of the form $x(z_0, z')$ where $z_0$ is an important latent semantic variable denoting object identity, while $z'$ is an unimportant latent variable denoting nuisance features, like object pose or location. The augmentation procedure scrambles
pose/position while preserving object identity. Consider a simple concrete example (Fig. 3(a)):

\[ \mathbf{x}(z_0, z') = \begin{cases} 
  \mathbf{e}_{z'} + \mathbf{e}_{(z'+1) \mod d} & z_0 = 1 \\
  \mathbf{e}_{z'} + \mathbf{e}_{(z'+2) \mod d} & z_0 = 2,
\end{cases} \]

(4)

Here \( 0 \leq z' \leq d - 1 \) denotes \( d \) discrete translational object positions on a periodic ring and \( z_0 \in \{1, 2\} \) denotes two possible objects \( 11 \) and \( 101 \). The distribution is uniform both over objects and positions: \( p(z_0, z') = \frac{1}{2d} \). Augmentation shifts the object to a uniformly random position via \( p_{\text{aug}}(z'|z_0) = 1/d \). For a single linear neuron \( K_1(\mathbf{x}) = \mathbf{x} \), and the augmentation-averaged connection is \( \overline{K}_1(z_0) = \frac{1}{2} \mathbf{1} \), and is actually independent of object identity \( z_0 \) (both objects activate two pixels at any location). Thus \( \mathbb{V}_{z_0} [\overline{K}_1(z_0)] = 0 \) and no learning happens.

A local receptive field (RF) does not help. In the same generative model, now consider a linear neuron with a local RF of width \( 2 \). Within the RF only four patterns can arise: \( 00, 01, 10, 11 \). Taking the expectation over \( z' \) given \( z_0 \) (Fig. 3(a2)) yields \( \overline{K}_1(z_0=1) = \frac{1}{d} [x_{11} + x_{01} + x_{10} + (d - 3)x_{00}] \) and \( \overline{K}_1(z_0=2) = \frac{1}{d} [2x_{01} + 2x_{10} + (d - 4)x_{00}] \). Here, \( x_{11} \in \mathbb{R}^2 \) denotes pattern \( 11 \). This yields

\[ \mathbb{V}_{z_0} [\overline{K}_1(z_0)] = \frac{1}{4d^2} \mathbf{u} \mathbf{u}^T \quad \text{where} \quad \mathbf{u} := x_{11} + x_{00} - x_{01} - x_{10}. \]

(5)

and \( \mathbb{V}_{z_0} [\overline{K}_1(z_0)] \in \mathbb{R}^{2 \times 2} \) since the RF has width \( 2 \). Note that the signed sum of the four pattern vectors in \( \mathbf{u} \) actually cancel, so that \( \mathbf{u} = 0 \), implying \( \mathbb{V}_{z_0} [\overline{K}_1(z_0)] = 0 \) and no learning happens. Interestingly, although the conditional distribution of the \( 4 \) input patterns depends on the object identity \( z_0 \) (Fig. 3(a2)), a linear neuron cannot learn to discriminate the objects.

A nonlinear neuron with local RF can learn to detect object selective features. With a ReLU neuron with weight vector \( \mathbf{w} \), from Def. 1, the connection is now \( \overline{K}_1(\mathbf{x}; \mathbf{w}) = \psi(\mathbf{w}^T \mathbf{x}) \mathbf{x} \). Suppose \( \mathbf{w}(t) \) happens to be selective for a single pattern \( \mathbf{x}_p \) (where \( p \in \{00, 01, 10, 11\} \)), i.e., \( \mathbf{w}(t)^T \mathbf{x}_p > 0 \) and \( \mathbf{w}(t)^T \mathbf{x}_{p'} < 0 \) for \( p' \neq p \). The augmentation averaged connection is then \( \overline{K}_1(z_0) \propto \mathbf{x}_p \) where the proportionality constant depends on object identity \( z_0 \) and can be read off (Fig. 3(a2)). Since this averaged connection varies with object identity \( z_0 \) for all \( p \), the covariance operator does not vanish and is given by \( \mathbb{V}_{z_0} [\overline{K}_1(z_0)] = \mathbb{V}_{z_0}[\mathbf{c}_p \mathbf{x}_p \mathbf{x}_p^T] \) where the constant \( \mathbb{V}_{z_0}[\mathbf{c}_p] = 0 \) depends on the selective pattern \( p \) and can be computed from Fig. 3(a2). By Theorem 3, the dot product \( \mathbf{x}_p^T \mathbf{w}(t) \) grows over time:

\[ \mathbf{x}_p^T \mathbf{w}(t + 1) = \mathbf{x}_p^T (I_{2 \times 2} + \alpha \beta \mathbf{c}_p \mathbf{x}_p \mathbf{x}_p^T) \mathbf{w}(t) = \left( 1 + \alpha \beta \mathbf{c}_p \mathbf{c}_p \mathbf{w}(t) > \mathbf{x}_p^T \mathbf{w}(t) \right). \]

(6)

Thus the learning dynamics amplifies the initial selectivity to the object selective feature vector \( \mathbf{x}_p \) in a way that cannot be done with a linear neuron. Note this argument also holds with bias terms and initial selectivity for more than one patterns. Moreover, with a local RF, the probability of weak initial selectivity to local object sensitive features is high, and we may expect amplification of such weak selectivity in real neural network training.

5.2 Two-layer case with multiple hidden neurons

Now consider a two-layer network \( (L = 2) \) with 1 output \( (n_L = 1) \) and 1-hidden layer with \( n_1 \) ReLU neurons (Fig. 2(c)). In this case, the augmentation-averaged connection \( \overline{K}_1(z_0) \) at the lowest layer \( l = 1 \) can be written as

\[ \overline{K}_1(z_0) = [w_{1,1} \mathbf{u}_1^T(z_0), w_{2,1} \mathbf{u}_1^T(z_0), \ldots, w_{2,n_1} \mathbf{u}_1^T(z_0)]^T \in \mathbb{R}^{n_1 \times d} \]

(7)

where \( \mathbf{u}_j(z_0) := \mathbb{E}_{z'|z_0} [\mathbf{x}(z_0, z') \mathbf{w}_j^T(z_0, z') \geq 0] \) is the expected activation of the hidden layer. Note that each \( \mathbf{u}_j(z_0) \in \mathbb{R}^d \) and \( w_{2,j} \) is the scalar weight of the second weight matrix that connects node \( j \) to the output. Then we have the following theorem:

**Theorem 4.** If \( \mathbb{Cov}_{z_0} [\mathbf{u}_j, \mathbf{u}_k] = 0 \) for \( j \neq k \), then the time derivative of \( w_{2,j} \) and \( \mathbf{w}_1, \mathbf{w}_2 \) satisfies:

\[ w_{2,j} = w_{2,j} \mathbf{A}_j \mathbf{w}_1, \quad \mathbf{w}_1 = w_{2,j} \mathbf{A}_j \mathbf{w}_1, \quad \text{where} \quad \mathbf{A}_j := \mathbb{V}_{z_0}[\mathbf{u}_j(z_0)]. \]

(8)
The set of all nodes and all latent variables at layer $l$. Nodes corresponding to latent variable $z_\mu, \mathcal{N}_\mu^{ch}$ are children under $\mathcal{N}_\mu$.

The top-down transition probability from $z_\mu$ to $z_\nu$.

Expected activation $f_j$ given $z_\mu$ ($z_\nu$’s descendants are marginalized).

Activations for all nodes $j \in \mathcal{N}_\mu$ and for the children of $\mathcal{N}_\mu$.

Polarity of the transitional probability.

Polarity of probability of root latent $z_0$.

Discrepancy of node $k$ w.r.t its latent variable $z_\nu(k)$.

Child selectivity vector.

Table 1: Notation for Sec. 5.3 (binary symmetric HLTM).

It is easy to see that $d|w_{2j}|^2/dt = d||w_{1j}||^2/dt$ and thus $w_{2j}^2 = ||w_{1j}||^2 + c$ where $c$ is some constant which doesn’t change over time. Since $A_j$ is PSD, $w_{1j}A_j w_{1j} \geq 0$ and $w_{2j}$ always grows, which in turn accelerates the convergence of $u_{1j}$ to the largest eigenvector of $A_j$. This shows top-down modulation: the existence of top-layer weights accelerates the training of the lower layer. Previous works (Allen-Zhu and Li, 2020) also mention a similar concept in supervised learning, called “backward feature corrections”, for deep polynomial networks with quadratic activations. Here we provide a rigorous formulation showing similar behaviors happen in SSL with SGD training.

Note that if we consider ReLU neurons, $A_j$ changes with $w_{1j}$; while for linear nodes, $A_j$ is a constant, since the gating $I(w_{1j}, x > 0)$ is always 1. If the data $x$ forms a mixture of two Gaussians: $x \sim \frac{1}{2} \mathcal{N}(z_0 = 1)N(w^*_1, \sigma^2) + \frac{1}{2} \mathcal{N}(z_0 = 0)N(w^*_2, \sigma^2)$ and let $\Delta w^* := w^*_1 - w^*_2$, then in the linear case, $A_j \sim \Delta w^* \Delta w^*^\top$ and $u_{1j}$ converges to $\pm \Delta w^*$ (Fig. 3(b)). In the nonlinear case with multiple Gaussians, if one of the Gaussians sits at the origin (e.g., background noise), then dependent on initialization, $A_j$ evolves into $w_{k}^* w_{k}^*\top$ for some center $k$, and $u_{1j} \rightarrow u_{k}^*$. Note this dynamics is insensitive to specific parametric forms of the input data.

5.3 DEEP RELU SSL TRAINING WITH HIERARCHICAL LATENT TREE MODELS (HLTM)

For multi-layer case, we mainly study a hierarchical latent tree model, where each leaf variable is sampled via a Markov Chain from the root latent variable $z_0$. The model is motivated by the structure of ConvNet: each latent variable at the intermediate layer generates a subset of the observable data $x$ that is sent to the deep ReLU networks, in which the nodes also have local receptive fields (Fig. 2(d)).

We define symbols in Tbl. 1. At layer $l$, we have categorical latent variables $\{z_l\}$, where $\mu \in Z_l$ is the Greek letter index for latent variables. Each $z_l$ can take discrete values. The topmost latent variable is $z_0$. Following the tree structure, for $\mu \in Z_l$ and $\nu_1, \nu_2 \in Z_{l-1}$, conditional independence holds: $P(z_{\nu_1}, z_{\nu_2} | z_\mu) = P(z_{\nu_1} | z_\mu)P(z_{\nu_2} | z_\mu)$. The final sample $x$ is the instantiation of the leaf latents (Fig. 2(d)), and thus depends on all latent variables. Corresponding to the hierarchical tree model, each neural network node $j \in N_l$ maps to a unique $\mu = \mu(j) \in Z_l$. Let $\mathcal{N}_\mu$ be all nodes that map to $\mu$. For $j \in \mathcal{N}_\mu$, its activation $f_j$ only depends on the value of $z_\mu$ and its descendant latent variables. We also define $v_j(z_\mu) := E_x[f_j | z_\mu]$ as the expected activation w.r.t $z_\mu$. Given a sample $x$, the data augmentation is done by resampling all $z_\mu$ (which are $z$ in Fig. 2) given the root $z_0$.

**Symmetric Binary HLTM**. Here we consider a symmetric binary case: each $z_l \in \{0, 1\}$ and for $\mu \in Z_l, \nu \in Z_{l-1}, P(z_\nu = 1 | z_\mu = 1) = P(z_\nu = 1 | z_\mu = 0) = (1 + \rho_{\mu\nu})/2$, where the polarity $\rho_{\mu\nu} \in [-1, 1]$ measures how informative $z_\nu = 1$ is. If $\rho_{\mu\nu} = \pm 1$ then there is no stochasticity in the top-down generation process; $\rho_{\mu\nu} = 0$ means no information in the downstream latents and the posterior of $z_\mu$ given the observation $x$ can only be uniform. See Appendix for more general cases.

Now we compute covariance operator $\nabla_{z_0}[\bar{K}_\mu(z_0)]$ at different layers, where $\bar{K}_\mu(z_0) = \mathbb{E}_{z_0'}[f_{\mathcal{N}_\mu^{ch}}(z_0)]$, we first check the term $\mathbb{E}_{z_0'}[f_{\mathcal{N}_\mu^{ch}}(z_0)]$ and incorporate the Jacobian later.

**Theorem 5** (Covariance operator in binary HLTM). $\nabla_{z_0}[\mathbb{E}_{z_0'}[f_{\mathcal{N}_\mu^{ch}}(z_0)]] = \frac{\rho_{\mu\nu}^2}{2}(1 - \rho_{\mu\nu}^2)a_{\mathcal{N}_\mu^{ch}}a_{\mathcal{N}_\mu^{ch}}^\top$.

Here $a_{\mathcal{N}_\mu^{ch}} := |\rho_{\mu\nu}(k), s_k|_{k \in N_\mu}$. If $\max_{\alpha, \beta} |\alpha, \beta| \leq \gamma < 1$, then $\lim_{L \to \infty} \rho_{\mu} \to 0$ for $\mu \in Z_l$.

Theorem 5 suggests when $\rho_{\mu\nu}$ and $||a_{\mathcal{N}_\mu^{ch}}||$ is large, the covariance operator is large and training is faster. For deep HLTM and deep networks, at lower layers, $\rho_{\mu\nu} \to 0$ and $P_{\theta_0}$ becomes uniform due to mixing from the progression of the Markov Chain, making $\nabla_{z_0}[\bar{K}_1(z_0)]$ small. Therefore, training in SSL is faster at the top layers where the covariance operators have large magnitude. On
the other hand, large $\|a_{N^\text{ch}_{\mu}}\|$ means that $s_k := (v_k(1) - v_k(0))/2$ is large, or the expected activation $v_k(z_\nu)$ is selective among different values of $z_\nu$ for $\nu \in \text{ch}(\mu)$. Interestingly, this can be achieved by over-parameterization (i.e., $|N_{\mu}| > 1$):

**Theorem 6** (Lucky nodes in deep ReLU networks regarding to binary HLTM). Suppose each element of the weights $W_l$ between layer $l + 1$ and $l$ are initialized with $\text{Uniform} \left[-\sigma_w \sqrt{3/|N^\text{ch}_{\mu}|}, \sigma_w \sqrt{3/|N^\text{ch}_{\mu}|}\right]$. There exists $\sigma_k^2$ so that $\forall f_k | z_\nu |^2 \leq \sigma_k^2$ for any $k \in N_l$.

For any $\mu \in \mathbb{Z}+1$, if $|N_{\mu}| = O(\exp(c))$, then with high probability, there exists at least one node $j \in N_{\mu}$ so that their pre-activation gap $|v_j(1) - v_j(0)| = 2w^j_1a_{N^\text{ch}_{\mu}} > 0$ and the activations satisfy:

$$
|v_j^2(1) - v_j^2(0)| \geq 3\sigma^2_w \left[ \frac{1}{4|N^\text{ch}_{\mu}|} \sum_{k \in N^\text{ch}_{\mu}} |v_k(1) - v_k(0)|^2 \left( \frac{e + 6}{6} \rho_{2\mu}^2 - 1 \right) - \sigma_k^2 \right].
$$

Intuitively, this means that with large polarity $\rho_{2\mu}$ (or strong signals sending top-down), over-parametrized ReLU networks with random initialization yields selective neurons, if the lower layer also contains selective ones. For example, when $\sigma_l = 0$, $c = 9$, if $\rho_{2\mu} \geq 63.3\%$ then there is a gap between expected activations $v_j(1)$ and $v_j(0)$, and the gap is larger when the selectivity in the lower layer $l$ is higher. Note that at the lowest layer, $\{v_k\}$ are themselves observable leaf latent variables and are selective by definition. So a bottom-up mathematical induction will happen.

After the gradient update, if we assume the top-down Jacobian is just 1, then for the “lucky” node $j$ we will have (for brevity, let $c := \alpha \beta \rho_{2\mu}^2 (1 - \rho_{2\mu}^2)$):

$$
a_{N^\text{ch}_{\mu}}^{T} w_j(t+1) = a_{N^\text{ch}_{\mu}}^{T} \left[ I + ca_{N^\text{ch}_{\mu}}a_{N^\text{ch}_{\mu}}^{T} \right] w_j(t) = (1 + c\|a_{N^\text{ch}_{\mu}}\|_2^2)a_{N^\text{ch}_{\mu}}^{T} w_j(t) > a_{N^\text{ch}_{\mu}}^{T} w_j(t) > 0
$$

which means that the pre-activation gap $|\tilde{v}_j(1) - \tilde{v}_j(0)| = 2w^j_1a_{N^\text{ch}_{\mu}}$ grows over time and the latent variable $z_\mu$ is learned (instantiated as $f_j$) during training, even if the network is never supervised with its true value. Similar to Sec. 5.2, once the top layer starts to have large weights, $K_1(z_0)$ for lower layer becomes larger due to large Jacobian and training is accelerated.

We implement the HLTM and confirm, as predicted by our theory, that the intermediate layers of deep ReLU networks do indeed learn the latent variables of the HLTM (see Tbl. 2 below).

### 6 Analysis of Ingredients underlying the Success of BYOL

In BYOL, the two networks are no-longer identical and, interestingly, only positive pairs are used for training. The first network with weights $W_1 = \mathcal{W} := \{W_\text{base}, W_\text{pred}\}$ is an online network that is trained to predict the output of the second target network with weights $W_2 = \mathcal{W}'$ using a learnable predictor $W_\text{pred}$ to map the online to target outputs (Fig. 1(a) and Fig. 1 in Grill et al. (2020)). In contrast, the target network has $W' := \{W_\text{base}, W'_\text{pred}\}$, where $W'_\text{pred}$ is an identity and $W_\text{base}$ is exponential moving average (EMA) of $W_\text{base}$: $W_\text{base}'(t+1) = \gamma_\text{ema} W_\text{base}'(t) + (1 -\gamma_\text{ema}) W_\text{base}(t)$.

We now theoretically analyze BYOL using the same theoretical framework. Since BYOL only uses positive pairs, we consider the following loss function for BYOL:

$$
r := \frac{1}{2} \|f_L(x_1; \mathcal{W}) - f_{L'}(x_+; \mathcal{W}')\|^2_2
$$

where the input data are positive pairs: $x_1, x_+ \sim p_{\text{amp}}(\cdot|\cdot)$ and $x \sim p(\cdot)$. The two outputs, $f_L(x_1; \mathcal{W})$ and $f_{L'}(x_+; \mathcal{W}')$, are from the online and the target networks, respectively. Note that $L' < L$ due to the presence of an extra predictor on the side of online network ($W$). Similar to SimCLR case, we leave the incorporation of $\ell_2$ normalization of the final layer as the future work.

Without EMA ($\gamma_\text{ema} = 0$) and the predictor, we obtain $\mathcal{W}' = \mathcal{W}$ and the BYOL update without BN for $\mathcal{W}$ is (derivation in Appendix E.1):

$$
\text{vec}(\Delta W_l)_{\text{sym}} = -\mathbb{E}_x \left[ \mathcal{W}_{x' \sim p_{\text{amp}}(\cdot|\cdot)} \left[ K_l(x') \right] \right] \text{vec}(W_l).
$$

This update only promotes variance minimization in the representations of different augmented views of the same data samples and therefore would yield model collapse.

We now consider the effects played by the extra predictor and BatchNorm (BN) (Ioffe and Szegedy, 2015) in BYOL. Our interest on BatchNorm is motivated by a recent blogpost (Fetterman and Albrecht, 2020). We will see that combining both could yield a sufficient condition to create an implicit
contrastive term for BYOL to work. As pointed recently by Rich mond et al. (2020), BYOL can still work using other normalization techniques that do not rely on cross batch statistics (e.g., GroupNorm (Wu and He, 2018), Weight Standardization (Qiao et al., 2019) and careful initialization of affine transform of activations). Here we demonstrate a “sufficient condition” using our framework, and leave analysis of other normalization techniques for future work.

In this paper, we analyze a simplified version of BN in which the mini-batch mean only is subtracted. When adding predictor, Theorem 1 can still be applied by adding identity layers on top of the target network \( \mathcal{W} \) so that the online and the target networks have the same depth. Theorem 5 in (Tian, 2018) demonstrates this version of BN shifts the downward gradients so their mini-batch mean is 0:

\[
\hat{g}_i := g_i - \frac{1}{|B|} \sum_{i \in B} g_i = g_i - \bar{g}_i
\]

Here \( g_i \) is the \( i \)-th sample in a batch and \( \bar{g}_i \) is the batch average (same for \( f_i \)). Backpropagating through this BN (vs. just subtracting the mean only in the forward pass\(^1\)), leads to a correction term:

**Theorem 7.** If (1) the network is linear from layer \( l \) to the topmost and (2) the downward gradient \( g_l \) undergoes Eqn. 12, then with large batch limits, the correction of the update is\(^2\) (for brevity, dependency on \( \mathcal{W} \) is omitted, while dependency on \( \mathcal{W}' \) is made explicit):

\[
\text{vec}(\delta W^l_{\text{BN}}) = \mathbb{E}_x [\bar{K}_i(x)] \{ \mathbb{E}_x [\bar{K}_i^T(x)] \text{vec}(W_i) - \mathbb{E}_x [\bar{K}_i^T(x; \mathcal{W})] \text{vec}(W'_i) \}
\]

and the corrected weight update is \( \Delta W_l := \Delta W_l + \delta W^l_{\text{BN}} \). Using Eqn. 11, we have:

\[
\text{vec}(\Delta W_l) = \text{vec}(\Delta W_{l_{\text{sym}}}) - \mathbb{V}_x [\bar{K}_i(x)] \text{vec}(W_i) + \text{Cov}_x [\bar{K}_i(x), \bar{K}_i(x; \mathcal{W})] \text{vec}(W'_i)
\]

Theorem 7 makes several predictions.

**SimCLR.** In this baseline case, both networks use the same weight and there is no predictor. This means \( \mathcal{W}' = \mathcal{W} \). From the analysis above, we have \( \delta W^l_{\text{BN}} = 0 \) for gradient through positive (and negative) pair distance \( r_+ \) (and \( r_- \)) and BN should not matter, modulo its benefit during optimization. This is justified in the recent blogpost (Fetterman and Albrecht, 2020).

**BYOL.** When the predictor is present, \( \mathcal{W}' \neq \mathcal{W} \) and BN is present, from the analysis above we know that \( \delta W^l_{\text{BN}} \neq 0 \), which provides an implicit contrastive term. Note that \( \mathcal{W}' \neq \mathcal{W} \) means there is a predictor, the online network uses EMA, or both.

**The Predictor.** We first discuss the predictor without EMA. To see why \( \mathcal{W}_{\text{pred}} \) plays a critical role, consider the last two terms (denoted as \( \bar{\Delta W}_l \)) in Eqn. 14:

\[
\bar{\Delta W}_l = -\mathbb{V}_x [\bar{K}_i(x; \mathcal{W})] \text{vec}(W_i) + \text{Cov}_x [\bar{K}_i(x), \bar{K}_i(x; \mathcal{W}), \bar{K}_i(x; \mathcal{W}')] \text{vec}(W'_i)
\]

If there is no EMA (i.e., \( \mathcal{W}'_{\text{base}} = \mathcal{W}_{\text{base}} \)) and the weights of the predictor are all small positive numbers (e.g., \( \approx \beta > 0 \)), then \( +\text{Cov}_x [\bar{K}_i(x), \bar{K}_i(x; \mathcal{W})] \approx \beta \mathbb{V}_x [\bar{K}_i(x; \mathcal{W}_{\text{base}})] \) and for all layer \( l \) in \( \mathcal{W}_{\text{base}} \), Eqn. 21 becomes:

\[
\bar{\Delta W}_l \approx \beta (1 - \beta) \mathbb{V}_x [\bar{K}_i(x; \mathcal{W}_{\text{base}})] \text{vec}(W_i)
\]

Intuitively, in Eqn. 21, the first term \( -\mathbb{V}_x [\bar{K}_i(x; \mathcal{W})] \) is second order in the Jacobian of \( \mathcal{W}_{\text{pred}} \), while the second term \( +\text{Cov}_x [\bar{K}_i(x; \mathcal{W}), \bar{K}_i(x; \mathcal{W}')] \) is first-order with respect to \( \mathcal{W}_{\text{pred}} \). So if the predictor weight \( \mathcal{W}_{\text{pred}} \) is “small” in magnitude (e.g., \( \beta \ll 1 \)), then the latter term dominates and \( \bar{\Delta W}_l \) becomes the covariance operator. In this regime, BYOL with predictor+BN sensibly amplifies variance across data samples through the covariance operator, and minimizes variance across different augmented views of the same data sample through \( -\mathbb{E}_x [\bar{K}_i(x; \mathcal{W})] \), which is the first term in Eqn. 14. Interestingly, SimCLR doesn’t have such a term.

**The Exponential Moving Average (EMA).** On the other hand, the EMA part might play a different role. Consider the following linear dynamic system, which is a simplified version of Eqn. 14:

\[
w(t+1) - w(t) = \Delta w(t) = \alpha [-w(t) + (1-\lambda)w_{\text{ema}}(t)]
\]

\(^1\)In PyTorch, the former is \( x-x.mean() \) and the latter is \( x-x.mean().detach() \).

\(^2\)A formal treatment requires Jacobian \( J \) to incorporate BatchNorm’s contribution and is left for future work.
Using z-transform (note here the symbol z has nothing to do with latent variables in Sec. 5), we could compute its two roots (See Appendix E.3 for the derivation):

$$z_{\min,\max} = 1 - \frac{1}{2} \left[ (1 - \gamma_{\text{ema}} + \alpha) \pm \sqrt{(1 - \gamma_{\text{ema}} + \alpha)^2 - 4\alpha(1 - \gamma_{\text{ema}})\lambda} \right]$$

As analyzed above, the magnification factor of $w_{\text{ema}}(t)$ is larger than that of $w(t)$, and thus $\lambda < 0$, and $z_{\max} > 1 > z_{\min}$. As a result, $w(t) \propto z_{\max}^t$ will have exponential growth and learning happens. Compared to no EMA case (i.e., $\gamma_{\text{ema}} = 0$), with a $\gamma_{\text{ema}} < 1$ but close to 1, $z_{\max}$ becomes smaller (but still > 1) and the exponential growth is less aggressive, which stabilizes the training.

Indeed, empirical findings in a recent blogpost (Fetterman and Albrecht, 2020) as well as our own experiments (Tbl. 3) suggests that standard BYOL without BN fails. In addition, we also initialize the predictor with small positive weights (See Appendix F), as well as reinitialize the predictor weight once in a while (Tbl. 5), and BYOL still works well. These empirical evidences are all consistent with the theoretical prediction.

Stop Gradient. In BYOL, the target network ($\mathcal{W}'$) serves as a target to be learned from, but does not contribute gradient to the current weight $\mathcal{W}$. While this seems to be obvious if we use EMA, when $\gamma_{\text{ema}} = 0$ (no EMA) and thus $\mathcal{W}'_{\text{base}} = \mathcal{W}_{\text{base}}$, intuitively it is not clear whether the target network should also contribute the gradient or not.

Fortunately, using our framework, we could clearly see that stop gradient is important. If we receive gradient from both the online and the target network, then for any common layer $l$, the weight update $\text{vec}(\Delta \bar{W}_l)$ becomes symmetric (note that this can be derived by swapping $\mathcal{W}'$ with $\mathcal{W}$ and add the two terms together):

$$\text{vec}(\Delta \bar{W}_l) = 2\text{vec}(\Delta W_l)_{\text{sym}} - \left( \nabla_{\mathcal{W}} \left[ \tilde{K}_l(x; \mathcal{W}) \right] + \nabla_{\mathcal{W}} \left[ \tilde{K}_l(x; \mathcal{W}') \right] \right) \text{vec}(W_l) + 2\text{Cov}_{\mathcal{W}} \left[ \tilde{K}_l(x; \mathcal{W}), \tilde{K}_l(x; \mathcal{W}') \right] \text{vec}(W_l)$$

In this case, since $\mathcal{W}$ contains the predictor and $\mathcal{W}'$ does not, the terms $-\nabla_{\mathcal{W}} \left[ \tilde{K}_l(x; \mathcal{W}) \right] - \nabla_{\mathcal{W}} \left[ \tilde{K}_l(x; \mathcal{W}') \right]$ now becomes zero-th order with respect to $\mathcal{W}_{\text{pred}}$ and dominates the weight update. On the other hand, the approximate covariance operator $+2\text{Cov}_{\mathcal{W}} \left[ \tilde{K}_l(x; \mathcal{W}), \tilde{K}_l(x; \mathcal{W}') \right]$ is first-order and is smaller. Therefore, the training will collapse and won’t learn anything.

Table 2: Normalized Correlation between the topmost latent variables in binary HLTM and topmost nodes in deep ReLU networks ($L = 5$) go up when training with SimCLR with NCE loss. We see higher correlations at both initialization and end of training with more over-parameterization (Left: $|\mathcal{W}| = 2$, Right: $|\mathcal{W}| = 5$).

| $\rho_{\mu,\sigma}$ | Initial | 1 epoch | 20 epochs | $\rho_{\mu,\sigma}$ | Initial | 1 epoch | 20 epochs |
|----------------------|---------|---------|-----------|----------------------|---------|---------|-----------|
| ~ Uniform[0.7, 1]    | 0.31    | 0.69    | 0.76      | ~ Uniform[0.7, 1]    | 0.60    | 0.72    | 0.88      |
| ~ Uniform[0.8, 1]    | 0.65    | 0.76    | 0.79      | ~ Uniform[0.8, 1]    | 0.73    | 0.80    | 0.87      |
| ~ Uniform[0.9, 1]    | 0.81    | 0.85    | 0.86      | ~ Uniform[0.9, 1]    | 0.87    | 0.90    | 0.95      |

Table 3: Top-1 STL performance with different combination of predictor (P), EMA and BatchNorm using BYOL. When EMA is on, we use $\gamma_{\text{ema}} = 0.996$. Batch size is 128 and all experiments run on 5 seeds.

| - | EMA | BN | EMA | BN | P | P | P | P | P | P | P | P |
|---|-----|----|-----|----|---|---|---|---|---|---|---|---|
| 38.7 ± 0.6 | 39.3 ± 0.9 | 33.0 ± 0.3 | 32.8 ± 0.5 | 39.5 ± 3.1 | 44.4 ± 3.2 | 63.6 ± 1.06 | 78.1 ± 0.3 |

Table 4: Top-1 STL performance using different BatchNorm components in the predictor and the projector of BYOL ($\gamma_{\text{ema}} = 0.996$, 100 epochs). There is no affine part. “$\mu, \sigma$” = BN without affine, “$\mu, \sigma^\square$” = normalization with mean and std but only backpropagating through mean. All variants with detached zero-mean normalization (in red) yield similar poor performance as no normalization.

| - | $\mu$ | $\sigma$ | $\mu, \sigma$ | $\mu^\square$ | $\sigma^\square$ | $\mu, \sigma^\square$ | $\mu^\square, \sigma^\square$ |
|---|-------|--------|---------------|-------------|-----------------|-----------------|------------------|
| 43.9 ± 4.2 | 64.8 ± 0.6 | 72.2 ± 0.9 | 78.1 ± 0.3 | 44.2 ± 7.0 | 54.2 ± 0.6 | 48.3 ± 2.7 | 76.3 ± 0.4 | 47.0 ± 8.1 |

7 EXPERIMENTS

We test our theoretical findings through experiments on STL-10 (Coates et al., 2011) and ImageNet (Deng et al., 2009). We use a simplified linear evaluation protocol: the linear classifier is
Table 5: Top-1 performance of BYOL using reinitialization of the predictor every $T$ epochs.

|                | Original BYOL | ReInit $T = 5$ | ReInit $T = 10$ | ReInit $T = 20$ |
|----------------|---------------|----------------|-----------------|-----------------|
| STL-10 (100 epochs) | 78.1          | 78.6           | 79.1            | 79.0            |
| ImageNet (60 epochs)  | 60.9          | 61.9           | 62.4            | 62.4            |

trained on frozen representations computed without data augmentation. This reuses pre-computed representations and accelerates evaluation by 10x.

**Hierarchical Latent Tree Model (HLTM).** We implement the HLTM and check whether the intermediate layers of deep ReLU networks learn the corresponding latent variables at the same layer. The degree of learning is measured by the normalized correlations between the ground truth latent variable $z_\mu$ and its best corresponding node $j \in N_\mu$ at each layer. Tbl. 2 indicates this measure increases with over-parameterization and learning, consistent with the analysis in Sec. 5.3.

**Factors underlying BYOL performance.** To test our theory, we perform an ablation study of BYOL on STL-10 by modifying three key components: predictor, EMA and BN. Tbl. 3 shows that BN and predictor are important and EMA further improves the performance. First, without a predictor, neither BN nor EMA give good performance. A predictor without BN still doesn’t work. A predictor with BN starts to show good performance (63.6%) and further adding EMA leads to the best performance (78.1%). This is consistent with our theoretical findings in Sec. 6, in which we show that using a predictor with BN yields $\delta W_l^{\text{BN}} \neq 0$ and leads to an implicit contrastive term.

To further test our understanding of the role played by BN, we fractionate BN into several sub-components: subtract by batch mean (mean-norm), divide by batch standard deviation (std-norm) and affine, and do ablation studies (Tbl. 4). Surprisingly, removing affine yields slightly better performance on STL-10 (from 78.1% to 78.7%). We also find that variants of mean-norm performs reasonably, while variants of detached mean-norm has similar poor performance as no normalization, supporting that centralizing backpropagated gradient leads to implicit contrastive terms (Sec. 6). Note that std-norm also helps, which we leave for future analysis.

We also check whether the online network requires an “optimal predictor” as suggested by recent version (v3) of BYOL. For this, we reinitialize the predictor (ReInit) every $T$ epochs and compare the final performance under linear evaluation protocol. Interestingly, as shown in Tbl. 5, ReInit actually improves the performance a bit, compared to the original BYOL that keeps training the same predictor, which should therefore be closer to optimal. Moreover, if we shrink the initial weight range of the predictor to make $\text{Cov}_x \left[ \bar{K}_l(x), \bar{K}_l(x; W') \right]$ (third term in Eqn. 14) more dominant, and reduce the learning rate, the performance further improves (See Tbl. 10 in Appendix F), thereby corroborating our analysis.

8 Conclusion and Future Works

In this paper, we propose a novel theoretical framework to study self-supervised learning (SSL) paradigm that consists of dual deep ReLU network. We analytically show that the weight update at each intermediate layer is govern by the covariance operator, a PSD matrix that amplifies the direction of initial weights that aligns with the difference across samples averaged by data augmentation. We show how the covariance operator interacts with multiple generative models that generate the input data distribution, including a simple 1D model with circular translation and hierarchical latent tree models. Finally, with the same framework, we give a theoretical explanation why BYOL with a predictor and BatchNorm can work without negative pairs. Experiments on both STL-10 and ImageNet support our theoretical findings.

To our best knowledge, our work is the first to bridge multiple critical components together, including contrastive loss, data augmentation, (hierarchical) generative models, self-supervision, deep non-linear models, regularization techniques (like BatchNorm) and emergence of features and representations. We hope this work opens new opportunities and perspectives for the research community.
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A BACKGROUND AND BASIC SETTING (SECTION 3)

A.1 THEOREM 1

Definition 2 (reversibility). A layer \( l \) is reversible if there is a \( G_l(x; W) \in \mathbb{R}^{n_l \times n_l} \) so that \( f_l(x; W) = G_l(x; W)f_{l-1}(x; W) \) and \( g_{l-1} = G_l^T(x; W)Q_l g_l \) for some constant PSD matrix \( \mathbb{R}^{n_l \times n_l} \ni Q_l \geq 0 \). A network is reversible if all layers are.

Note that many different kinds of layers have this reversible property, including linear layers (MLP and Conv) and (leaky) ReLU nonlinearity. For multi-layer ReLU network, for each layer \( l \), we have:

\[
G_l(x; W) = D_l(x; W)W_l, \quad Q_l \equiv I_{n_l \times n_l}
\]

(22) where \( D_l \in \mathbb{R}^{n_l \times n_l} \) is a binary diagonal matrix that encodes the gating of each neuron at layer \( l \). The gating \( D_l(x; W) \) depends on the current input \( x \) and current weight \( W \).

In addition to ReLU, other activation function also satisfies this condition, including linear, LeakyReLU and monomial activations. For example, for power activation \( \psi(x) = x^p \) where \( p > 1 \), we have (where \( f_l \) is the pre-activation at layer \( l \)):

\[
G_l(x; W) = \text{diag}^{p-1}(\tilde{f}_l)W_l, \quad Q_l \equiv pI_{n_l \times n_l}
\]

(23)

Remark. Note that the reversibility is not the same as invertible. Specifically, reversibility only requires the transfer function of a backpropagation gradient is a transpose of the forward function.

Lemma 1 (Recursive Gradient Update (Extension to Lemma 1 in (Tian, 2020)). Let (pseudo)-Jacobian matrix \( J_l(x) = I_{n_l \times n_l} \) and recursively define \( J_{l-1}(x) := J_l(x)\sqrt{Q_l}G_l(x) \in \mathbb{R}^{n_l \times n_l} \). Here \( \sqrt{Q_l} \) is the constant PSD matrix so that \( \sqrt{Q_l}^T\sqrt{Q_l} = Q_l \geq 0 \).

If (1) the network is reversible (Def. 2) and (2) \( \sqrt{Q_l} \) commutes with \( J_l(x_1)\top J_l(x_1) \) and \( J_l(x_1)\top J_l(x_2) \), then minimizing the \( \ell_2 \) objective

\[
r(W_l) := \frac{1}{2} \| f_L(x_1; W_l) - f_L(x_2; W_l) \|^2
\]

(24) with respect to weight matrix \( W_l \) at layer \( l \) yields the following gradient at layer \( l \):

\[
g_l = J_l^T(x_1) [J_l(x_1) f_l(x_1) - J_l(x_2) f_l(x_2)]
\]

(25)

Proof. We prove by induction. Note that our definition of \( W_l \) is the transpose of \( W_l \) defined in (Tian, 2020). Also our \( g_l(x) \) is the gradient before nonlinearity, while (Tian, 2020) uses the same symbol for the gradient after nonlinearity.

For notation brievity, we let \( f_l(x_1) := f_l(x_1; W_l) \) and \( G_l(x_1) := G_l(x_1; W_l) \). Similar for \( x_2 \) and \( W_2 \).

When \( l = L \), by the property of \( \ell_2 \)-loss, we know that \( g_L = f_L(x_1; W_l) - f_L(x_2; W_2) \), by setting \( J_L(x_1) = J_L(x_2) = I \), the condition holds. Now suppose for layer \( l \), we have:

\[
g_l = J_l^T(x_1) [J_l(x_1) f_l(x_1) - J_l(x_2) f_l(x_2)]
\]

(26)

Then:

\[
g_{l-1} = G_l^T(x_1)Q_l g_l
\]

(27)

\[
= G_l^T(x_1)Q_l J_{l-1}(x_1) [J_l(x_1) f_l(x_1) - J_l(x_2) f_l(x_2)]
\]

(28)

\[
= G_l^T(x_1)\sqrt{Q_l} J_{l-1}(x_1) \left[ J_l(x_1)\sqrt{Q_l} f_l(x_1) - J_l(x_2)\sqrt{Q_l} f_l(x_2) \right]
\]

(29)

\[
= J_{l-1}^T(x_1) \left[ J_{l-1}(x_1)\sqrt{Q_l} G_l(x_1) f_{l-1}(x_1) - J_{l-1}(x_2)\sqrt{Q_l} G_l(x_2) f_{l-1}(x_2) \right]
\]

(30)

\[
= J_{l-1}^T(x_1) \left[ J_{l-1}(x_1) f_{l-1}(x_1) - J_{l-1}(x_2) f_{l-1}(x_2) \right]
\]

(31)

Note that for multi-layered ReLU network, \( G_l(x) = D_l(x)W_l, Q_l = I \) for each ReLU+Linear layer, if we set \( x_1 = x_2 = x, W_1 = W, W_2 = W^* \) (teacher weights), then we go back to the original Lemma 1 in (Tian, 2020). \( \square \)
Remark on the top-most $\ell_2$-normalization layer. For $\ell_2$-normalized layer $f_l := f_{l-1}/\|f_{l-1}\|_2$, we have $G_l := 1/\|f_{l-1}\|_2 \cdot I_{n_l \times n_l}$ and due to the following identity (here $\hat{y} := y/\|y\|_2$):

$$
\frac{\partial \hat{y}}{\partial y} = \frac{1}{\|y\|_2^2} (I - \hat{y}\hat{y}^T)
$$

(32)

Therefore we have $\partial f_l / \partial f_{l-1} = (I_{n_l \times n_l} - f_l f_l^T) G_l$ and we could set $Q_l := I - f_l f_l^T$, which is a projector matrix and thus PSD. Furthermore, since the normalization layer is at the topmost, $J_l = I$ and $Q_l$ trivially commutes with $J_l^T J_l$.

The only issue is that $Q_l$ is not a constant PSD matrix and can change over training. Therefore Lemma 1 doesn’t apply exactly to such a layer but can be regarded as an approximate way to model.

Remark on ResNet. Note that the same structure holds for blocks of ResNet with ReLU activation.

Now we prove Theorem 1 in a more general setting where the network is reversible (note that deep ReLU networks are included and its $Q_l$ is a simple identity matrix):

**Lemma 2** (Squared $\ell_2$ Gradient for dual deep reversible networks). The gradient $g_{W_l}$ of the squared loss $r$ with respect to $W_l \in \mathbb{R}^{n_l \times n_{l-1}}$ for a single input pair $(x_1, x_2)$ is:

$$
g_{W_l} = \text{vec}(\partial r / \partial W_{l,1}) = K_{1,l} \left[ K_{1,l}^T \text{vec}(W_{l,1}) - K_{2,l}^T \text{vec}(W_{l,2}) \right].
$$

(33)

Here $K_l(x; W) := f_{l-1}(x; W) \otimes J_l^T(x; W)$. $K_{1,l} := K_l(x_1; W_{l,1})$ and $K_{2,l} := K_l(x_2; W_{l,2})$.

**Proof.** We consider more general case where the two towers have different parameters, namely $W_1$ and $W_2$. Applying Lemma 1 for the branch with input $x_1$ at the linear layer $l$, and we have:

$$
g_{1,l} = J_l^T [J_{1,l} W_{l,1} f_{l-1} - J_{2,l} W_{l,2} f_{l-1}]
$$

(34)

where $f_{l-1} := f_{l-1}(x_1; W_{l,1})$ is the activation of layer $l - 1$ just below the linear layer at tower 1 (similar for other symbols). The gradient (and the weight update, according to gradient descent) of the weight $W_l$ between layer $l$ and layer $l - 1$ is:

$$
\frac{\partial r}{\partial W_{l,1}} = g_{1,l} J_{1,l}^T
$$

(35)

$$
= J_{1,l}^T J_{1,l} W_{l,1} f_{l-1} - J_{1,l}^T J_{2,l} W_{l,2} f_{l-1} - J_{1,l}^T J_{1,l} W_{l,2} f_{l-1} + J_{1,l}^T J_{2,l} W_{l,1} f_{l-1}
$$

(36)

Using $\text{vec}(AXB) = (B^T \otimes A) \text{vec}(X)$ (where $\otimes$ is the Kronecker product), we have:

$$
\text{vec} \left( \frac{\partial r}{\partial W_{l,1}} \right) = \left( f_{l-1} J_{1,l}^T - f_{l-1} J_{2,l}^T \right) \text{vec}(W_{l,1}) - \left( f_{l-1} J_{1,l}^T - f_{l-1} J_{2,l}^T \right) \text{vec}(W_{l,2})
$$

(37)

Let

$$
K_l(x; W) := f_{l-1}(x; W) \otimes J_l^T(x; W) \in \mathbb{R}^{n_l \times n_{l-1} \times n_{l-1}}
$$

(38)

Note that $K_l(x; W)$ is a function of the current weight $W$, which includes weights at all layers. By the mixed-product property of Kronecker product $(A \otimes B)(C \otimes D) = AC \otimes BD$, we have:

$$
\text{vec} \left( \frac{\partial r}{\partial W_{l,1}} \right) = K_l(x_1) K_l(x_1)^T \text{vec}(W_{l,1}) - K_l(x_1) K_l(x_2)^T \text{vec}(W_{l,2})
$$

(39)

$$
= K_l(x_1) [K_l(x_1)^T \text{vec}(W_{l,1}) - K_l(x_2)^T \text{vec}(W_{l,2})]
$$

(40)

where $K_l(x_1) = K_l(x_1; W_{l,1})$ and $K_l(x_2) = K_l(x_2; W_{l,2})$.

In SimCLR case, we have $W_1 = W_2 = W$ so

$$
\text{vec} \left( \frac{\partial r}{\partial W_l} \right) = K_l(x_1) [K_l(x_1) - K_l(x_2)]^T \text{vec}(W_l)
$$

(41)

\[\square\]
B Analysis of SimCLR using Teacher-Student Setting (Section 4)

B.1 Theorem 2

Proof. We have:
\[
\frac{\partial L}{\partial r_+} = \frac{1}{\tau} \left( 1 - \frac{e^{-r_+ / \tau}}{e^{-r_+ / \tau} + \sum_{k'=1}^H e^{-r_{k'} / \tau}} \right) > 0 \tag{42}
\]
\[
\frac{\partial L}{\partial r_{k-}} = - \frac{1}{\tau} \left( \frac{e^{-r_{k-} / \tau}}{e^{-r_+ / \tau} + \sum_{k'=1}^H e^{-r_{k'} / \tau}} \right) < 0, \quad k = 1, \ldots, H \tag{43}
\]
and obviously we have:
\[
\frac{\partial L}{\partial r_+} + \sum_{k=1}^H \frac{\partial L}{\partial r_{k-}} = 0 \tag{44}
\]
\[\square\]

B.2 Theorem 3

Proof. First we have:
\[
\text{vec}(g_{W_1}) = \frac{\partial L}{\partial W_1} = \frac{\partial L}{\partial r_+} + \sum_{k=1}^H \frac{\partial L}{\partial r_{k-}} \tag{45}
\]
Then we compute each terms. Using Theorem 1, we know that:
\[
\frac{\partial r_+}{\partial W_1} = K_1(x_1)(K_1(x_1) - K_1(x_+))^\top \text{vec}(W_1) \tag{46}
\]
\[
\frac{\partial r_{k-}}{\partial W_1} = K_1(x_1)(K_1(x_1) - K_1(x_{k-}))^\top \text{vec}(W_1), \quad k = 1, \ldots, n \tag{47}
\]
Since Eqn. 44 holds, $K_1(x_1)K_1^\top(x_1)$ will be cancelled out and we have:
\[
\text{vec}(g_{W_1}) = K_1(x_1) \sum_{k=1}^H \left[ \frac{\partial L}{\partial r_{k-}} (K_1(x_1) - K_1(x_{k-}))^\top \right] \text{vec}(W_1) \tag{48}
\]
Then we use the assumption that $\frac{\partial L}{\partial r_{k-}} = -\beta / H$ with $\beta > 0$:
\[
\text{vec}(g_{W_1}) = -\beta K_1(x_1) \left[ K_1^\top(x_+) - \frac{1}{n} \sum_{k=1}^H K_1^\top(x_{k-}) \right] \text{vec}(W_1) \tag{49}
\]
Taking large batch limits, we know that $\mathbb{E}[K_1(x_1)K_1^\top(x_+)] = \mathbb{E}_x[K_1(x)K_1^\top(x)]$ since $x_1, x_+ \sim p_{aug}(x)$ are all augmented data points from a common sample $x$. On the other hand, $\mathbb{E}[K_1(x_1)K_1^\top(x_{k-})] = \mathbb{E}_x[K_1(x)] \mathbb{E}_x[K_1^\top(x)]$ since $x_1$ and $x_{k-}$ are generated from independent samples and data augmentation. Therefore,
\[
\text{vec}(g_{W_1}) = -\beta \mathbb{E}_x \left[ \mathbb{E}_x[K_1(x)K_1^\top(x)] - \mathbb{E}_x[K_1(x)] \mathbb{E}_x[K_1^\top(x)] \right] \text{vec}(W_1) \tag{50}
\]
\[
= -\beta \mathbb{E}_x \left[ K_1(x) \right] \text{vec}(W_1) \tag{51}
\]
\[\square\]

C The Dynamics of Two-layer ReLU Network and the Interplays of Covariance Operators between Nearby Layers (Section 5.2)

C.1 Theorem 4

Proof. For convenience, we define the centralized version of $u_j(z_0)$: $\hat{u}_j(z_0) = u_j(z_0) - \mathbb{E}_{z_0}[u_j(z_0)]$ and the matrices $A_{jk} := \text{Cov}_{z_0}[u_j(z_0), u_k(z_0)] = \mathbb{E}_{z_0}[\hat{u}_j(z_0)\hat{u}_k^\top(z_0)]$. 

17
At layer $l = 1$ the covariance operator is $\mathbb{V}_{z_0} [\hat{K}_1(z_0)] = [w_{2,j} w_{2,k} A_{jk}] \in \mathbb{R}^{n_1 d \times n_1 d}$.

On the other hand, if we check the second layer $l = 2$, we could compute $\hat{K}_2(z_0) \in \mathbb{R}^{n_1}$:

$$
\hat{K}_2(z_0) = \mathbb{E}_{z' \mid z_0} [f_1(z_0) = \mathbb{E}_{z' \mid z_0} [D_1 W_1 x \mid z_0] = \mathbb{E}_{z' \mid z_0} [x^\top \otimes D_1(x) \mid z_0] \text{vec}(W_1)
$$

(52)

For each component $j$, we have $[\hat{K}_2(z_0)]_j = w_{1,j}^T u_j(z_0)$, or:

$$
\hat{K}_2(z_0) = [w_{1,1}^T u_1(z_0), w_{1,2}^T u_2(z_0), \ldots, w_{1,n_1}^T u_{n_1}(z_0)]
$$

(53)

So at layer $l = 2$ we can compute the covariance operator $\mathbb{V}_{z_0}[\hat{K}_2(z_0)] = [w_{1,j}^T A_{jk} w_{1,k}] \in \mathbb{R}^{n_1 \times n_1}$.

Using the assumption that $A_{jk} := \mathbb{E}_{z_0} [u_j(z_0) u_k^*(z_0)] = \mathbb{E}_{z_0} [u_j(z_0)] \mathbb{E}_{z_0} [u_k^*(z_0)] = 0$ for $j \neq k$. We arrive at the conclusion that $w_j$ is only coupled with $v_j$, yielding the following dynamical system:

$$
\dot{w}_{2,j} = w_{2,j} w_{1,j}^T A_{jj} w_{1,j}, \quad \dot{w}_{1,j} = w_{2,j}^2 A_{jj} w_{1,j}
$$

(54)

### D Hierarchical Latent Tree Models (Section 5.3)

#### D.1 Lemmas

**Lemma 3** (Variance Squashing). Suppose a function $\phi : \mathbb{R} \mapsto \mathbb{R}$ is $L$-Lipschitz continuous: $|\phi(x) - \phi(y)| \leq L|x - y|$, then for $x \sim p(\cdot)$, we have:

$$
\mathbb{V}_p[\phi(x)] \leq L^2 \mathbb{V}_p[x]
$$

(55)

**Proof.** Suppose $x, y \sim p(\cdot)$ are independent samples and $\mu_\phi := \mathbb{E} [\phi(x)]$. Note that $\mathbb{V} [\phi(x)]$ can be written as the following:

$$
\mathbb{E} [\phi(x) - \phi(y)]^2 = \frac{1}{2} \mathbb{E} [(\phi(x) - \mu_\phi) - (\phi(y) - \mu_\phi)]^2
$$

$$
= \mathbb{E} [\phi(x) - \mu_\phi]^2 + \mathbb{E} [(\phi(y) - \mu_\phi)^2] - 2 \mathbb{E} [(\phi(x) - \mu_\phi)(\phi(y) - \mu_\phi)]
$$

$$
= 2 \mathbb{V}_p[\phi(x)]
$$

Therefore we have:

$$
\mathbb{V}_p[\phi(x)] = \frac{1}{2} \mathbb{E} [\phi(x) - \phi(y)]^2 \leq \frac{L^2}{2} \mathbb{E} [|x - y|^2] = L^2 \mathbb{V}_p[x]
$$

(57)

**Lemma 4** (Sharpened Jensen’s inequality (Liao and Berg, 2018)). If function $\phi$ is twice differentiable, and $x \sim p(\cdot)$, then we have:

$$
\frac{1}{2} \mathbb{V}_p[x] \inf \phi'' \leq \mathbb{E} [\phi(x)] - \phi(\mathbb{E} [x]) \leq \frac{1}{2} \mathbb{V}_p[x] \sup \phi''
$$

(58)

**Lemma 5** (Sharpened Jensen’s inequality for ReLU activation). For ReLU activation $\psi(x) := \max(x, 0)$ and $x \sim p(\cdot)$, we have:

$$
0 \leq \mathbb{E} [\psi(x)] - \psi(\mathbb{E} [x]) \leq \sqrt{\mathbb{V}_p[x]}
$$

(59)

**Proof.** Since $\psi$ is a convex function, by Jensen’s inequality we have $\mathbb{E} [\psi(x)] - \psi(\mathbb{E} [x]) \geq 0$. For the other side, let $\mu := \mathbb{E} p [x]$ and we have (note that for ReLU, $\psi(x) - \psi(\mu) \leq |x - \mu|$):

$$
\mathbb{E} [\psi(x)] - \psi(\mathbb{E} [x]) = \int (\psi(x) - \psi(\mu)) p(x) dx
$$

(60)

$$
\leq \int |x - \mu| p(x) dx
$$

(61)

$$
\leq \left( \int |x - \mu|^2 p(x) dx \right)^{1/2} \left( \int p(x) dx \right)^{1/2}
$$

(62)

$$
= \sqrt{\mathbb{V}_p[x]}
$$

(63)

where the last inequality is due to Cauchy-Schwarz.

\[\square\]
D.2 MORE GENERAL ASSUMPTION OF CONDITIONAL DISTRIBUTION IN HLT

We consider a more general case of HLT where each \( z_\mu \) is categorical: \( 0 \leq z_\mu < m_\mu \). For convenience, we define the following symbols for \( k \in \mathcal{N}_{\mu}^{ch} \) (note that \( |\mathcal{N}_{\mu}^{ch}| = N_{\mu}^{ch} \) is the number of the children of the node set \( \mathcal{N}_{\mu} \)):

\[
\begin{align*}
\nu_{\mu k} & := \mathbb{E}_z [f_j(z_\mu)] = P_{\mu \nu}(k) \nu_k \in \mathbb{R}^{m_\mu} \\
V_{\mu, \mathcal{N}_{\mu}^{ch}} & := [\nu_{\mu k}]_{k \in \mathcal{N}_{\mu}^{ch}} \\
\tilde{v}_j & := \left[ \mathbb{E}_z \{f_j(z_\mu)\} \right] = V_{\mu, \mathcal{N}_{\mu}^{ch}} w_j \in \mathbb{R}^{m_\mu}.
\end{align*}
\]

As an extension of binary symmetric HLT, we make an assumption for the transitional probability:

**Assumption 1.** For \( \mu \in \mathcal{Z}_1 \) and \( \nu \in \mathcal{Z}_{l-1} \), the transitional probability matrix \( P_{\mu \nu} := [P(z_\nu | z_\mu)] \) has decomposition \( P_{\mu \nu} = \frac{1}{m_\nu} 1_\mu 1_\nu^\top + C_{\mu \nu} \) where \( C_{\mu \nu} 1_\nu = 0_\mu \) and \( 1_\mu^\top C_{\mu \nu} = 0_\nu \).

Note that \( C_{\mu \nu} 1_\nu = 0_\mu \) is obvious due to the property of conditional probability. The real condition is \( 1_\mu^\top C_{\mu \nu} = 0_\nu \). If \( m_\mu = m_\nu \), then \( P_{\mu \nu} \) is a square matrix and Assumption 1 is equivalent to \( P_{\mu \nu} \) is double-stochastic. Assumption 1 makes computation of \( P_{\mu \nu} \) easy for any \( z_\mu \) and \( z_\nu \).

**Lemma 6** (Transition Probability). If Assumption 1 holds, then for \( \mu \in \mathcal{Z}_1, \nu \in \mathcal{Z}_{l-1} \) and \( \alpha \in \mathcal{Z}_{l-2} \), we have:

\[
P_{\mu \alpha} = P_{\mu \nu} P_{\nu \alpha} = \frac{1}{m_\alpha} 1_\mu 1_\alpha^\top + C_{\mu \alpha} C_{\nu \alpha}
\]

In general, for any \( \mu \in \mathcal{N}_{l_1} \) and \( \alpha \in \mathcal{N}_{l_2} \) with \( l_1 > l_2 \), we have:

\[
P_{\mu \alpha} = \frac{1}{m_\alpha} 1_\mu 1_\alpha^\top + \prod_{\mu, \ldots, \zeta, \ldots, \alpha} C_{\zeta \zeta}
\]

**Proof.** Using Assumption 1, we have

\[
P_{\mu \alpha} = P_{\mu \nu} P_{\nu \alpha}
\]

\[
= \left( \frac{1}{m_\nu} 1_\mu 1_\nu^\top + C_{\mu \nu} \right) \left( \frac{1}{m_\alpha} 1_\nu 1_\alpha^\top + C_{\nu \alpha} \right)
\]

since \( 1_\nu 1_\nu = m_\nu, C_{\mu \nu} 1_\nu = 0_\mu \) and \( 1_\nu^\top C_{\nu \alpha} = 0_\alpha \), the conclusion follows.

**Remark.** In the symmetric binary HLT mentioned in the main text, all \( C_{\mu \nu} \) can be parameterized as (here \( q := [1, 1]^T \)):

\[
C_{\mu \nu} = C_{\mu \nu}(\rho_{\mu \nu}) = \frac{1}{2} \begin{bmatrix} \rho_{\mu \nu} & -\rho_{\mu \nu} \\ -\rho_{\mu \nu} & \rho_{\mu \nu} \end{bmatrix} = \frac{1}{2} \rho_{\mu \nu} qq^T
\]
Moreover, since $1^T_2 C_{\mu \nu} = 0_2$ and $C_{\mu \nu} 1_2 = 0_2$ provides 4 linear constraints (1 redundant), leaving 1 free parameter, which is the polarity $\rho_{\mu \nu} \in [-1,1]$ of latent variable $z_\nu$ given its parent $z_\mu$. Moreover, since $q^T q = 2$, the parameterization is close under multiplication:

$$C(\rho_{\mu \nu})C(\rho_{\nu \alpha}) = \frac{1}{4} q q^T q q^T \rho_{\mu \nu} \rho_{\nu \alpha} = \frac{1}{2} q q^T \rho_{\mu \nu} \rho_{\nu \alpha} = C(\rho_{\mu \nu} \rho_{\nu \alpha})$$

(72)

D.3 Theorem 5

Proof. First note that for each node $k \in N_{\mu}^\circ$: 

$$v_{ik} := [\mathbb{E}_z [f_k | z_0]] = \sum_{z_\nu} \mathbb{E}_z [f_k | z_\nu] \mathbb{P}(z_\nu | z_0) = P_0 v_k \tag{73}$$

$$= \left( \frac{1}{m_\nu} 1_1^T + C_{0\nu} \right) v_k \tag{74}$$

$$= \frac{1}{m_\nu} 1_1^T v_k + C_{0\nu} v_k \tag{75}$$

Note that $1_1^T v_k$ is a constant regarding to change of $z_0$. So we could remove it when computing covariance operator. On the other hand, for a categorical distribution 

$$v_{ik} := [\mathbb{E}_z [f_k | z_0]] = \sum_{z_\nu} \mathbb{E}_z [f_k | z_\nu] \mathbb{P}(z_\nu | z_0) = P_0 v_k \tag{73}$$

$$= \left( \frac{1}{m_\nu} 1_1^T + C_{0\nu} \right) v_k \tag{74}$$

$$= \frac{1}{m_\nu} 1_1^T v_k + C_{0\nu} v_k \tag{75}$$

With $P_0 := \text{diag}(\mathbb{P}(z_0))$, the mean is $\mathbb{E}_z_0 [u] = 1^T P_0 u$ and its covariance can be written as (here $1 = 1_1$):

$$\mathbb{V}_0 [u_k] = (u - 11^T P_0 u) P_0 (u - 11^T P_0 u) = u^T (P_0 - P_0 11^T P_0) u \tag{76}$$

Note that each column of $V_{0,N_\mu^{ch}}$ is $v_{0k}$. Setting $u = v_{0k}$ and we have:

$$\mathbb{V}_0 [\mathbb{E}_z [f_{N_\mu^{ch}} | z_0]] = V_{0,N_\mu^{ch}} (P_0 - P_0 11^T P_0) V_{0,N_\mu^{ch}} \tag{77}$$

Note that $\mathbb{E}_z_0 [v_{0k}] = \frac{1}{m_\nu} 1_1^T v_k + \mathbb{E}_z_0 [C_{0\nu} v_k]$, since $11^T P_0 1 = 1$. With some computation, we could see $\text{Cov}_{z_0} [v_{0k}, v_{0k}'] = \text{Cov}_{z_0} [C_{0\nu} v_k, C_{0\nu} (k') v_{k'}]$. The equation above can be applied for any cardinality of latent variables. In the binary symmetric case, we have (note here we define $\rho_0 := \mathbb{P}(z_0 = 1) - \mathbb{P}(z_0 = 0)$, and $q := [-1,1]^T$):

$$P_0 - P_0 11^T P_0 = \frac{1}{4} (1 - \rho_0^2) q q^T \tag{78}$$

Note that in the binary symmetric case, according to remarks in Lemma 6, all $C_{\mu \nu} = \frac{1}{2} \rho_{\mu \nu} q q^T$ and we could compute $C_{0\nu} v_k$:

$$C_{0\nu} v_k = \frac{1}{2} \rho_{0\nu} q q^T v_k = \rho_{0\nu} \frac{1}{2} (v_k(1) - v_k(0)) q = \rho_{0\nu} s_k q \tag{79}$$

where according to Eqn. 72, we have:

$$\rho_{0\nu} := \prod_{\alpha, \beta, \ldots, \nu} \rho_{\alpha \beta} \tag{80}$$

and the covariance between node $k$ and $k'$ can be computed as:

$$\text{Cov}_{z_0} [v_{0k}, v_{0k'}] = \text{Cov}_{z_0} [C_{0\nu} v_k, C_{0\nu} (k') v_{k'}] \tag{81}$$

$$= \rho_{0\nu(k)} \rho_{0\nu(k')} s_k s_{k'} \frac{1}{4} q q^T q (1 - \rho_0^2) \tag{82}$$

$$= \rho_{0\nu(k)} \rho_{0\nu(k')} s_k s_{k'} (1 - \rho_0^2) \tag{83}$$

$$= \rho_{0\nu(k)} \rho_{0\nu(k')} s_k s_{k'} (1 - \rho_0^2) \tag{84}$$

The last equality is due to the fact that due to tree structure, the path from $z_0$ to all child nodes in $N_\mu^{ch}$ must pass $z_\mu$.

Therefore we can compute the covariance operator:

$$\mathbb{V}_0 [\mathbb{E}_z [f_{N_\mu^{ch}} | z_0]] = \rho_{0\nu}^2 (1 - \rho_0^2) a_{N_\mu^{ch}} a_{N_\mu^{ch}}^T \tag{85}$$

When $L \to +\infty$, we have:

$$\rho_{0\nu} := \prod_{\alpha, \beta, \ldots, \nu} \rho_{\alpha \beta} \to 0 \tag{86}$$

and thus the covariance becomes zero as well. \qed
D.4 Theorem 6

Proof. According to our setting, for each node $k \in \mathcal{N}_\mu$, there exists a unique latent variable $z_\nu$ with $\nu = \nu(k)$ that corresponds to it. In the following we omit its dependency on $k$ for brevity.

Since we are dealing with binary case, we define the following for convenience:

\begin{align}
v_k^+ &:= v_k(1) \\
v_k^- &:= v_k(0) \\
\bar{v}_k &:= \frac{1}{2}(v_k^+ + v_k^-) = \frac{1}{2}(\mathbb{E}_z[f_k|z_\nu = 1] + \mathbb{E}_z[f_k|z_\nu = 0]) \\
\bar{v}^{\bar{v}}_{\mathcal{N}_\mu} &:= [\bar{v}_k]_{k \in \mathcal{N}_\mu} \in \mathbb{R}^{\mathcal{N}_\mu} \\
v_k &:= \frac{1}{2} (v_k^+ - v_k^-) = \frac{1}{2}(\mathbb{E}_z[f_k|z_\nu = 1] - \mathbb{E}_z[f_k|z_\nu = 0]) \\
\bar{s}^{\bar{s}}_{\mathcal{N}_\mu} &:= [s_k]_{k \in \mathcal{N}_\mu} \in \mathbb{R}^{\mathcal{N}_\mu}
\end{align}

We also define the sensitivity of node $k$ to be $\lambda_k := |(v_k^+)^2 - (v_k^-)^2|$. Intuitively, a large $\lambda_k$ means that the node $k$ is sensitive for changes of latent variable $z_\nu$. If $\lambda_k = 0$, then the node $k$ is invariant to latent variable $z_\nu$.

We first consider pre-activation $\tilde{f}_j := \sum_k w_{jk} f_k$ and its expectation with respect to latent variable $z$:

\[
\bar{v}^+_j := \mathbb{E}_z[\tilde{f}_j|z_\mu = 1], \quad \bar{v}^-_j := \mathbb{E}_z[\tilde{f}_j|z_\mu = 0]
\]

Note that for each node $k \in \mathcal{N}_\mu$ we have:

\[
v^+_{jk} = \bar{v}_k + \rho_{\mu} s_k, \quad v^-_{jk} = \bar{v}_k - \rho_{\mu} s_k
\]

Let $a^{\bar{a}}_{\mathcal{N}_\mu} := [a_k]_{k \in \mathcal{N}_\mu} := [\rho_{\mu} s_k]_{k \in \mathcal{N}_\mu}$ and

\[
a^{\bar{a}}_{\mathcal{N}_\mu} := V^{\bar{a}}_{\mu, \mathcal{N}_\mu} := [\mathbb{E}[f_k|z_\mu = 1]] = \bar{v}^{\bar{a}}_{\mathcal{N}_\mu} + a^{\bar{a}}_{\mathcal{N}_\mu}
\]

\[
a^{\bar{a}}_{\mathcal{N}_\mu} := V^{\bar{a}}_{\mu, \mathcal{N}_\mu} := [\mathbb{E}[f_k|z_\mu = 0]] = \bar{v}^{\bar{a}}_{\mathcal{N}_\mu} - a^{\bar{a}}_{\mathcal{N}_\mu}
\]

Then we have $\bar{v}^+_j = u^+_j u^{\bar{a}}_{\mathcal{N}_\mu}$ and $\bar{v}^-_j = u^+_j u^{\bar{a}}_{\mathcal{N}_\mu}$.

Note that $\bar{u}_j$ is a random variable with each entry $w_{jk} \sim \text{Uniform} \left[ -3\sigma_w \sqrt{3/|\mathcal{N}_\mu|}, 3\sigma_w \sqrt{3/|\mathcal{N}_\mu|} \right]$.

It is easy to verify that $\mathbb{E}[w_{jk}] = 0$ and $\text{Var}[w_{jk}] = \sigma_w^2/|\mathcal{N}_\mu|$. Therefore, for two dimensional vector $\bar{v}_j = [\bar{v}^+_j, \bar{v}^-_j]^T$, we can compute its first and second order moments: $\mathbb{E}[\bar{v}_j] = 0$ and $\text{Var}[\bar{v}_j] = \sigma^2/|\mathcal{N}_\mu|$.

Define the positive and negative set (note that $a_k := \rho_{\mu} s_k$):

\[
A_+ = \{k : a_k \geq 0\}, \quad A_- = \{k : a_k < 0\}
\]

Without loss of generality, assume that $\sum_{k \in A_+} a_k^2 \geq \sum_{k \in A_-} a_k^2$. In the following, we show there exists $j$ with $\lambda_j$ is greater than some positive threshold. Otherwise the proof is symmetric and we can show $\lambda_j$ is lower than some negative threshold.

When $|\mathcal{N}_\mu|$ is large, by Central Limit Theorem, $\nu$ can be regarded as zero-mean 2D Gaussian distribution and we have for some $c > 0$:

\[
P\left( \bar{v}^+_j \geq \frac{\sqrt{c} \sigma_w}{\sqrt{|\mathcal{N}_\mu|}} \|u^{\bar{a}}_{\mathcal{N}_\mu}\| \right) = \frac{1 - \text{erf}(\sqrt{c}/2)}{2}
\]

Moreover, if $a_l \neq 0$, then the following probability is also not small:

\[
P\left( \bar{v}^+_j \geq \frac{\sqrt{c} \sigma_w}{\sqrt{|\mathcal{N}_\mu|}} \|u^{\bar{a}}_{\mathcal{N}_\mu}\| \text{ and } \bar{v}^-_j < 0 \right)
\]
Therefore, when $|N_\mu| = O(\exp(c))$, with high probability, there exists $w_j$ so that

$$\tilde{v}_j^+ = w_j^T u_{\mu +} \geq \frac{\sqrt{c\sigma_w}}{\sqrt{|N_{\mu +}|}} |u_{\mu +}|, \quad \tilde{v}_j^- = w_j^T u_{\mu -} < 0$$  \hspace{1cm} (100)$$

Since $\tilde{v}_{\mu +} \geq 0$ (all $f_k$ are after ReLU and non-negative), this leads to:

$$\tilde{v}_j^+ \geq \frac{\sqrt{c\sigma_w}}{\sqrt{|N_{\mu +}|}} |u_{\mu +}| \geq \sqrt{c\sigma_w} \sum_{k \in A_v} a_k^2 \geq \sigma_w \frac{c}{2|N_{\mu +}|} \sum_{k \in N_{\mu +}} \rho_{\mu v}^2 s_k^2$$  \hspace{1cm} (101)$$

By Jensen’s inequality, we have (note that $\psi(x) := \max(x, 0)$ is the ReLU activation):

$$v_j^+ = E_z [f_j \mid z_\mu = 1] = E_z \left[ \psi(\tilde{f}_j) \mid z_\mu = 1 \right]$$  \hspace{1cm} (102)$$

$$\geq \psi \left( E_z \left[ \tilde{f}_j \mid z_\mu = 1 \right] \right) = \psi(\tilde{v}_j^+) \geq \sigma_w \frac{c}{2|N_{\mu +}|} \sum_{k \in N_{\mu +}} \rho_{\mu v}^2 s_k^2$$  \hspace{1cm} (103)$$

On the other hand, we also want to compute $v_j^- = E_z [f_j \mid z_\mu = 0]$ using sharpened Jensen’s inequality (Lemma 5). For this we need to compute the conditional covariance $V_z[\tilde{f}_j \mid z_\mu]$:

$$V_z[\tilde{f}_j \mid z_\mu] \equiv \sum_k w_{jk}^2 V_z[f_k \mid z_\mu] \leq \frac{3\sigma_w^2}{|N_{\mu +}|} \sum_k V_z[f_k \mid z_\mu]$$  \hspace{1cm} (104)$$

$$= \frac{3\sigma_w^2}{|N_{\mu +}|} \sum_k (E_{z_\mu \mid z_\nu} [V[f_k \mid z_\nu]] + V_{z_\mu \mid z_\nu} [E_z [f_k \mid z_\nu]])$$  \hspace{1cm} (105)$$

$$\leq \frac{3\sigma_w^2}{|N_{\mu +}|} \left( \sigma_f^2 + \frac{1}{|N_{\mu +}|} \sum_k V_{z_\mu \mid z_\nu} [E_z [f_k \mid z_\nu]] \right)$$  \hspace{1cm} (106)$$

Note that (104) is due to conditional independence: $f_k$ as the computed activation, only depends on latent variable $z_\nu$ and its descendants. Given $z_\mu$, all $z_\nu$ and their respective descendants are independent of each other and so does $f_k$. (105) is due to the fact that each $w_{jk}$ is sampled from uniform distribution and $|w_{jk}| \leq \sigma_w \frac{3}{|N_{\mu +}|}$.

Here $V_{z_\mu \mid z_\nu} [E_z [f_k \mid z_\nu]] = s_k^2 (1 - \rho_{\mu v}^2)$ can be computed analytically. It is the variance of a binomial distribution: with probability $\frac{1}{2} (1 + \rho_{\mu v})$ we get $v_k^+$ otherwise get $v_k^-$. Therefore, we finally have:

$$V_z[\tilde{f}_j \mid z_\mu] \leq 3\sigma_w^2 \left( \sigma_f^2 + \frac{1}{|N_{\mu +}|} \sum_k s_k^2 (1 - \rho_{\mu v}^2) \right)$$  \hspace{1cm} (107)$$

As a side note, using Lemma 3, since ReLU function $\psi$ has Lipschitz constant $\leq 1$ (empirically it is smaller), we know that:

$$V_z[f_j \mid z_\mu] \leq 3\sigma_w^2 \left( \sigma_f^2 + \frac{1}{|N_{\mu +}|} \sum_k s_k^2 (1 - \rho_{\mu v}^2) \right)$$  \hspace{1cm} (108)$$

Finally using Lemma 5 and $\tilde{v}_j^- < 0$, we have:

$$v_j^- = E_z [f_j \mid z_\mu = 0] = E_z \left[ \psi(\tilde{f}_j) \mid z_\mu = 0 \right]$$  \hspace{1cm} (109)$$

$$\leq \psi \left( E_z \left[ \tilde{f}_j \mid z_\mu = 0 \right] \right) + \sqrt{V_z[\tilde{f}_j \mid z_\mu = 0]}$$  \hspace{1cm} (110)$$

$$= \sqrt{V_z[\tilde{f}_j \mid z_\mu = 0]}$$  \hspace{1cm} (111)$$

$$\leq \sigma_w \sqrt{3\sigma_f^2 + \frac{3}{|N_{\mu +}|} \sum_k s_k^2 (1 - \rho_{\mu v}^2)}$$  \hspace{1cm} (112)$$

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Combining Eqn. 103 and Eqn. 112, we have a bound for $\lambda_j$:

$$\lambda_j = (v_j^+) - (v_j^-) \geq 3s_w^2 \left[ \frac{1}{N_{\text{ch}}} \sum_k s_k \left( \frac{c + 6}{6} \rho_{\mu \nu}^2 - 1 \right) \right]$$  \hfill (113)

\[\square\]

E THE ANALYSIS OF BYOL IN SEC. 6

E.1 DERIVATION OF BYOL GRADIENT

Note that for BYOL, we have:

$$\text{vec} \left( \frac{\partial r}{\partial W_l} \right) = K_l(x; W)[K_l^T(x; W)\text{vec}(W_l) - K_l^T(x)\text{vec}(W_l')]$$ \hfill (114)

under large batchsize, we have (note that we omit $W$ for any term that depends on $W$, but make dependence of $W'$ explicit in the math expression):

$$\text{vec} \left( \frac{\partial r}{\partial W_l} \right) = \mathbb{E}_{x \sim p(x)} \mathbb{E}_{x' \sim p_{aug}(\cdot|x)} [K_l(x')K_l^T(x') \text{vec}(W_l) - K_l(x)K_l^T(x; W')\text{vec}(W_l')]$$

For brevity, we write $\mathbb{E}_x [\cdot] := \mathbb{E}_{x \sim p(x)} [\cdot]$ and $\mathbb{E}_{x'} [\cdot] := \mathbb{E}_{x' \sim p_{aug}(\cdot|x)} [\cdot]$. Similar for $V$. And the equation above can be written as:

$$\text{vec} \left( \frac{\partial r}{\partial W_l} \right) = \mathbb{E}_x [V_x [K_l(x)]] \text{vec}(W_l) \hfill (115)$$

$$+ \mathbb{E}_x \{ K_l(x) [K_l^T(x)\text{vec}(W_l) - K_l^T(x; W')\text{vec}(W_l')] \} \hfill (116)$$

In terms of weight update by gradient descent, since $\Delta W_l = -\frac{\partial r}{\partial W_l}$, we have:

$$\text{vec} (\Delta W_l) = -\mathbb{E}_x [V_x [K_l(x)]] \text{vec}(W_l) \hfill (117)$$

$$- \mathbb{E}_x \{ K_l(x) [K_l^T(x)\text{vec}(W_l) - K_l^T(x; W')\text{vec}(W_l')] \} \hfill (118)$$

If we consider the special case $W = W'$, then the last two terms cancelled out, yielding:

$$\text{vec}(\Delta W_l)_{\text{sym}} = -\mathbb{E}_x [V_x [K_l(x)]] \text{vec}(W_l) \hfill (119)$$

And the general update (Eqn. 121) can be written as:

$$\text{vec} (\Delta W_l) = \text{vec}(\Delta W_l)_{\text{sym}} \hfill (120)$$

$$- \mathbb{E}_x \{ K_l(x) [K_l^T(x)\text{vec}(W_l) - K_l^T(x; W')\text{vec}(W_l')] \} \hfill (121)$$

E.2 THEOREM 7

Proof. When BN is present, Eqn. 114 needs to be corrected with an additional term, $\frac{\partial r}{\partial W_l} := \frac{\partial r}{\partial W_l} - \delta W_l^{\text{BN}}$, where $\delta W_l^{\text{BN}}$ is defined as follows:

$$\delta W_l^{\text{BN}} := \frac{1}{|B|} \sum_{i \in B} D_i^T g_i f_i^T$$  \hfill (122)

From the proof of Theorem 1, we know that for each sample $i \in B$:

$$D_i^T g_i = J_i^T [J_i W_i f_i^T - J_i (W') W_i f_i^T (W')]$$ \hfill (123)

Since the network is linear from layer $l$ to the topmost layer $L$, we have $D_i = \tilde{D}_i$. Since the only input dependent part in $J_i$ is the gating function between the current layer $l$ and the topmost layer $L$, for linear network the gating is always 1 and thus $\tilde{J}_i = J_i$ and is independent of input data. We
now have (note that we omit \( W \) for any terms that are dependent on \( W \), but will write \( W' \) explicitly for terms that are dependent on \( W' \)):

\[
\delta W_{l, BN}^{\text{BN}} := \frac{1}{|B|} \sum_{i \in B} D_l^i \delta g_l^i f_l^i \tau_{i-1} = -\bar{D}_l \delta g_l f_l \tau_{l-1} \tag{124}
\]

\[
\delta W_{l, BN}^{\text{BN}} = J_l^T (\bar{J}_l W_l f_{l-1} - \bar{J}_l (W') W_l' f_{l-1} (W')) f_{l-1, \tau_{l-1}} \tag{125}
\]

Therefore we have:

\[
\text{vec}(\delta W_{l, BN}^{\text{BN}}) = (\bar{f}_{l-1} \otimes J_l^T) \left[ (\bar{f}_{l-1} \otimes J_l^T) \text{vec}(W_l) - (\bar{f}_{l-1} (W') \otimes J_l^T (W')) \text{vec}(W_l') \right] \tag{126}
\]

Note that by assumption, since \( \bar{J}_l \) doesn’t depend on the input data, we have:

\[
\bar{f}_{l-1} \otimes J_l^T = E_B [f_{l-1}] \otimes J_l^T = E_B [f_{l-1} \otimes J_l^T] \tag{127}
\]

Taking large batchsize limits and notice that the batch \( B \) could contain any augmented data generated from independent samples from \( p(\cdot) \), we have:

\[
\text{vec}(\delta W_{l, BN}^{\text{BN}}) = E_{x, x'} [K_l(x')] E_{x, x'} [K_l^T (x')] \text{vec}(W_l) - E_{x, x'} [K_l(x)] E_{x, x'} [K_l^T (x'; W')] \text{vec}(W_l') \tag{128}
\]

\[
- E_{x, x'} [K_l(x')] E_{x, x'} [K_l^T (x'; W')] \text{vec}(W_l') \tag{129}
\]

An important thing is that the expectation is taking over \( x \sim p(x) \) and \( x' \sim p_{\text{aug}}(\cdot|x) \). Intuitively, this is because \( f_{l-1} \) and \( g_l \) are averages over the entire batch, which has both intra-sample and inter-sample variation.

With augment-mean connection \( \bar{K}_l(x) \) we could write:

\[
\text{vec}(\delta W_{l, BN}^{\text{BN}}) = E_x [\bar{K}_l(x)] E_x [\bar{K}_l^T (x')] \text{vec}(W_l) - E_x [\bar{K}_l(x)] E_x [\bar{K}_l^T (x; W')] \text{vec}(W_l') \tag{130}
\]

Plug in \( \delta W_{l, BN}^{\text{BN}} \) into Eqn. 121 and we have corrected gradient for BYOL:

\[
\text{vec} \left( \frac{\partial r}{\partial W_l} \right) = \text{vec} \left( \frac{\partial r}{\partial W_l} \right) - \text{vec}(\delta W_{l, BN}^{\text{BN}}) \tag{131}
\]

\[
= E_x \left[ \nabla_{x' \sim p_{\text{aug}}(\cdot|x)} [K_l(x')] \right] \text{vec}(W_l) + E_x \left[ \bar{K}_l(x) \right] \text{vec}(W_l) \tag{132}
\]

\[
- \text{Cov}_x [\bar{K}_l(x), \bar{K}_l(x'; W')] \text{vec}(W_l') \tag{133}
\]

And the weight update \( \Delta \tilde{W}_l = \Delta W_l + \delta W_{l, BN}^{\text{BN}} \) is:

\[
\text{vec}(\Delta W_l) = -E_x \left[ \nabla_{x' \sim p_{\text{aug}}(\cdot|x)} [K_l(x')] \right] \text{vec}(W_l) - E_x [\bar{K}_l(x)] \text{vec}(W_l) \tag{134}
\]

\[
+ \text{Cov}_x [\bar{K}_l(x), \bar{K}_l(x'; W')] \text{vec}(W_l') \tag{135}
\]

Using Eqn. 119, we have:

\[
\text{vec}(\Delta \tilde{W}_l) = \text{vec}(\Delta W_l) + \text{vec}(\delta W_{l, BN}^{\text{BN}}) \tag{136}
\]

\[
= \text{vec}(\Delta W_l)_{\text{sym}} \tag{137}
\]

\[
- \nabla_x [\bar{K}_l(x)] \text{vec}(W_l) + \text{Cov}_x [\bar{K}_l(x), \bar{K}_l(x'; W')] \text{vec}(W_l') \tag{138}
\]

\[
\square
\]

### E.3 Analysis on EMA

Consider the following discrete dynamics of a weight vector \( w(t) \):

\[
w(t+1) - w(t) = \alpha [-w(t) + (1 - \lambda)w_{ema}(t)] \tag{139}
\]

where \( \alpha \) is the learning rate, \( w_{ema}(t+1) = \gamma_{ema} w_{ema}(t) + (1 - \gamma_{ema}) w(t) \) is the exponential moving average of \( w(t) \). For convenience, we use \( \eta := 1 - \gamma_{ema} \).
Since it is a recurrence equation, we apply \( z \)-transform on the temporal domain, where \( w(z) := Z[w(t)] = \sum_{t=0}^{\infty} w(t) z^{-t} \). This leads to:

\[
 z(w(z) - w(0)) = w(z) - \alpha (w(z) - Z[w_{ema}(t)](1 - \lambda))
\]  

(140)

Note that for \( w_{ema}(t) \) we have:

\[
 z(w_{ema}(z) - w_{ema}(0)) = (1 - \eta)w_{ema}(z) + \eta w(z)
\]  

(141)

If we set \( w_{ema}(0) = 0 \), i.e., the target network is all zero at the beginning, then it gives \( w_{ema}(z) = \frac{\eta}{z - 1 + \eta} w(z) \). Plugging it back to Eqn. 140 and we have:

\[
 z(w(z) - w(0)) = w(z) - \alpha w(z) \left( 1 - \frac{\eta}{z - 1 + \eta}(1 - \lambda) \right)
\]  

(142)

And then we could solve \( w(z) \):

\[
 w(z) = \frac{z(z - 1 + \eta)}{(z - 1)^2 + (\eta + \alpha)(z - 1) + \alpha \eta \lambda} w(0)
\]  

(143)

Note that the denominator has two roots \( z_1 \) and \( z_2 \):

\[
 z_{1,2} = 1 - \frac{1}{2} \left( \eta + \alpha \pm \sqrt{\eta + \alpha)^2 - 4 \alpha \eta \lambda} \right)
\]  

(144)

and \( w(z) \) can be written as

\[
 w(z) = \frac{z(z - 1 + \eta)}{(z - z_1)(z - z_2)} w(0)
\]  

(145)

Without loss of generality, let \( z_1 < z_2 \). The larger root \( z_2 > 1 \) when \( \lambda < 0 \), so the zero \( z = 1 - \eta = \gamma_{ema} \) in the numerator won’t cancel out the pole at \( z_2 \). And we have:

\[
 \frac{z}{(z - z_1)(z - z_2)} = \frac{z}{z_2 - z_1} \left( \frac{z - z_1}{z - z_2} \right) - \frac{z - z_2}{z_2 - z_1} \left( \frac{1}{z - z_1} \right)
\]  

(146)

\[
 = \frac{z}{z_2 - z_1} \left( \frac{1}{1 - z_2 z^{-1}} - \frac{1}{1 - z_1 z^{-1}} \right)
\]  

(147)

where \( 1/(1 - z_2 z^{-1}) \) corresponds to a power series \( z_2 \) in the temporal domain. Therefore, we could see \( w(t) \) has exponential growth due to \( z_2 > 1 \).

**F ADDITIONAL EXPERIMENTS**

Based on our theoretical analysis, we try training the predictor in different ways and check whether it still works.

From the analysis in Sec. 6, we know that the reason why BYOL works is due to the dominance of \( \text{Cov}_x[K_t(x), \tilde{K}_t(x; W)] \) and its resemblance of the covariance operator \( \mathbb{V}_x[K_t(x)] \), which is a PSD matrix.

The dominance should be stronger if the predictor has smaller weights than normally initialized using Xavier/Kaiming initialization. Also, \( \text{Cov}_x[K_t(x), \tilde{K}_t(x; W')] \) should behave more like a PSD matrix, if the predictor’s weights are all small positive numbers and no BN is used.

The following table justifies our theoretical findings. In particular, Tbl. 10 shows better performance in STL-10 with smaller learning rate and smaller sample range of the predictor weights.
Table 7: Training one-layer predictor with positive initial weights and no EMA ($\gamma_{ema} = 0$). All experiments run for 3 seeds.

| Sample range of predictor weight | [0, 0.01] | [0, 0.02] | [0, 0.05] |
|----------------------------------|-----------|-----------|-----------|
| With BN in predictor             | 62.78 ± 1.40 | 62.94 ± 1.03 | 62.31 ± 1.80 |
| Without BN in predictor          | 71.95 ± 0.27 | 72.06 ± 0.44 | 71.91 ± 0.59 |

Table 8: Training one-layer predictor with positive initial weights with EMA ($\gamma_{ema} = 0.996$) and predictor resetting every $T = 10$ epochs. All experiments run for 3 seeds. Note that Xavier range is Uniform[$-0.15, 0.15$] and our initialization range is much smaller than that.

| Sample range of predictor weight | [0, 0.003] | [0, 0.005] | [0, 0.007] |
|----------------------------------|-----------|-----------|-----------|
| With BN in predictor             | 65.61 ± 1.34 | 70.56 ± 0.57 | 70.87 ± 1.51 |
| Without BN in predictor          | 74.39 ± 0.67 | 74.52 ± 0.63 | 74.80 ± 0.57 |

Table 9: Same as Tbl. 8 but with different weight range. All experiments run for 3 seeds.

| Sample range of predictor weight | [0, 0.01] | [0, 0.02] | [0, 0.05] |
|----------------------------------|-----------|-----------|-----------|
| With BN in predictor             | 68.98 ± 2.34 | 66.56 ± 1.70 | 68.41 ± 1.19 |
| Without BN in predictor          | 74.66 ± 0.81 | 73.60 ± 0.32 | 74.34 ± 0.77 |

Table 10: Top-1 Performance on STL-10 with a two-layer predictor with BN and EMA ($\gamma_{ema} = 0.996$). Learning rate is smaller (0.02) and predictor weight sampled from Uniform[$-\text{range}, \text{range}$]. Note that for this, Xavier range is Uniform[$-0.097, 0.097$] and our range is smaller.

| Weight range | $0.01$ | $0.02$ | $0.03$ | $0.05$ |
|--------------|-------|-------|-------|-------|
| $T = 3$      | 79.48 ± 0.40 | 79.70 ± 0.47 | 79.66 ± 0.37 | 78.63 ± 0.10 |
| $T = 5$      | 78.97 ± 0.62 | 79.63 ± 0.23 | 79.65 ± 0.37 | 79.01 ± 0.27 |
| $T = 10$     | 79.25 ± 0.20 | 79.63 ± 0.22 | 79.58 ± 0.25 | 79.18 ± 0.22 |
| $T = 20$     | 79.15 ± 0.66 | 79.91 ± 0.10 | 79.78 ± 0.05 | 79.54 ± 0.25 |