# On the Global Convergence of Natural Actor-Critic with Neural Network Parametrization 

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#### Abstract

Despite the empirical effectiveness of natural actor-critic (NAC) algorithms, their theoretical underpinnings remain relatively unexplored, especially with neural network parameterizations. In the existing literature, the non-asymptotic sample complexity bounds for NAC hold only when the critic is either tabular or are represented by a linear function. In this work, we relax such assumptions for NAC and utilize multi-layer neural network parameterization of the critic and an arbitrary smooth function for the actor. We establish the non-asymptotic sample complexity bounds of $\tilde{\mathcal{O}}\left(\frac{1}{\epsilon^{4}(1-\gamma)^{4}}\right)$ for the global convergence of NAC algorithm. We obtain this result using our unique decomposition of the error incurred at each critic step. The critic error is decomposed into the error incurred in fitting the sampled data, the error incurred due to the lack of knowledge of the transition matrix as well as the error incurred due to the limited approximation power of the class of neural networks. In contrast to the existing works for NAC with neural network parameterization of the critic, our analysis does not require i.i.d sampling.


## 1 InTRODUCTION

The use of neural networks in actor-critic (AC) algorithms is widespread in various machine learning applications, such as games (Vinyals et al., 2017; Bonjour et al., 2022), robotics (Morgan et al., 2021), autonomous driving (Kiran et al., 2022), ride-sharing (Al-Abbasi et al., 2019), networking (Geng et al., 2020), and recommender systems (Li et al., 2020). AC algorithms sequentially update the estimate of the actor (policy) and the critic (value function) based on the data collected at each iteration, as described in Konda \& Tsitsiklis (1999). An empirical and theoretical improvement of the AC, known as natural actor-critic (NAC), was proposed in Peters \& Schaal (2008). NAC replaced the stochastic gradient step of the actor with the natural gradient descent step described in Kakade (2001b) based on the theory of Rattray et al. (1998). The finite-time or non-asymptotic sample complexity bounds for NAC are limited to settings the critic has a linear parametrization (Xu et al., 2020a). However, since linear parametrization is quite restrictive, in practice mostly non-linear neural network-based parameterizations for both actor and critic are used as in Wang et al. (2021a). Despite the widespread use in practice, no finite-time sample complexity bounds are available for the setting when neural networks (NNs) represent the critic in the NAC algorithm.
The linear parametrization allows for a closed-form update for both the actor and the critic update steps. On the other hand, no such closed-form expressions are available for the NAC algorithm with non-linear parametrization. In recent work by Fu et al. (2020), a finite time bound is derived for linear function approximation but for NN parametrization, only an asymptotic convergence is established, moreover, this work requires i.i.d sampling. Wang et al. (2019) establishes similar asymptotic bounds where both the actor and critic are represented using a 2-layer neural network. Hence, we ask this question

## Is it possible to obtain non-asymptotic sample complexity bounds for global convergence of the natural actor-critic algorithm with a multi-layer neural network parametrization of the critic?

We answer this question by deriving precise non-asymptotic sample complexity bounds for the global convergence of the NAC algorithm. Our approach relies on decomposing the error incurred

Table 1: This table summarizes the sample complexities of different natural actor-critic algorithms. Our result is the first to provide sample complexity results of NAC for a general MDP setting with neural network (NN) parametrization for the critic.

| References | Actor <br> parametrization | Critic <br> parametrization | Sample <br> Complexity |
| :---: | :---: | :---: | :---: |
| (Xu et al., 2020b) | Linear | Linear | $\tilde{\mathcal{O}}\left(\epsilon^{-4}(1-\gamma)^{-9}\right)$ |
| (Khodadadian et al., 2021) | Linear | Linear | $\tilde{\mathcal{O}}\left(\epsilon^{-3}(1-\gamma)^{-11}\right)$ |
| (Xu et al., 2020a) | Linear | Linear | $\tilde{\mathcal{O}}\left(\epsilon^{-2}(1-\gamma)^{-4}\right)$ |
| (Wang et al., 2019) | 2-layer NN | 2-layer NN | Asymptotic |
| (Fu et al., 2020) | Multi-layer NN | Multi-layer NN | Asymptotic |
| This work | Multi-layer NN | Multi-layer NN | $\tilde{\mathcal{O}}\left(\epsilon^{-4}(1-\gamma)^{-4}\right)$ |

at each step of the NAC algorithm into the errors at the actor and critic steps separately. The main novelty in our approach is that the error incurred in the critic step is decomposed into the error in fitting the observed data, the error incurred due to the lack of knowledge of the transition matrix and the error due to finite approximation power of the class of neural networks. This contrasts the approach in Fu et al. (2020) where both the critic and actor optimizations are analyzed as stochastic gradient descent problems; thus, only an asymptotic error bound is possible with their approach. Hence, we summarize our contributions as follows.

- We derive a non-asymptotic sample complexity bound of $\tilde{\mathcal{O}}\left(\epsilon^{-4}(1-\gamma)^{-4}\right)$ for the global convergence of the natural actor-critic algorithm with neural network parameterizations for the critic and the actor. To achieve that, our two main novelties in the convergence analysis are highlighted next.
- Building upon the insights presented in Agarwal et al. (2021), we leverage the inherent smoothness property of the actor parametrization to derive an upper bound on the estimation error of the optimal value function. This upper bound is expressed in terms of the error incurred in attaining the compatible function approximation term, as elucidated in Sutton et al. (1999) and the error incurred in estimating the action value function used to solve the compatible function approximation.
- The error incurred at the critic step in fitting the data obtained through sampling is upper bounded using results from Allen-Zhu et al. (2019). The error incurred due to the lack of knowledge of the transition matrix is bounded in terms of the Radamacher complexity of the class of neural networks. This approach allows us to achieve the first non-asymptotic sample complexity bound for NAC with the critic parameterised by a multi layer neural network. It also allows us to have a milder assumption on the error incurred due to the limited approximation the function class representing the critic as compared to other finite time convergence results such as Xu et al. (2020a). Finally, we do not need to assume i.i.d sampling in this approach.


## 2 Related Works

Natural Policy Gradient. The problem of the non-convexity of the critic can be avoided if we use the natural policy gradient algorithm (Kakade, 2001b) where instead of maintaining a parameterized estimate of the critic, we obtain an estimate (or multiple estimates) at each iteration through a Monte Carlo estimate. In such a case, sample complexity estimates are possible without the assumption of linear function approximation of the value function. Agarwal et al. (2021) obtained a sample complexity bound of $\tilde{\mathcal{O}}\left(\frac{1}{\epsilon^{4}(1-\gamma)^{8}}\right)$, which was improved to $\tilde{\mathcal{O}}\left(\frac{1}{\epsilon^{2}(1-\gamma)^{7}}\right)$ in (Yuan et al., 2022) with the restriction of the actor being represented by a log-linear class of functions. Further improvement was obtained in (Liu et al., 2020b) with a sample complexity of $\tilde{\mathcal{O}}\left(\frac{1}{\epsilon^{3}(1-\gamma)^{6}}\right)$ and also did not require the restriction to log-linear class of functions to represent the actor. In spite of obtaining finite time sample complexity bounds, the Natural Policy Gradient algorithms suffer from high variance due to the Monte Carlo estimate. Additionally, each estimate of the critic requires on average a sample of size $\left(\frac{1}{1-\gamma}\right)$, thus these algorithms are not sample efficient in terms of $\gamma$. Additionally the
error incurred due to the Monte Carlo sampling of the critic as well as the lack of expressability of the class of functions representing the policy is represented as a constant.

Actor-Critic Methods. First conceptualized in Sutton (1988), aim to combine the benefits of the policy gradient methods and $Q$-learning based methods. The policy gradient step in these methods is replaced by a Natural Policy Gradient proposed in (Kakade, 2001b) to obtain the so-called Natural Actor Critic in (Peters et al., 2005). Sample complexity results for Actor Critic were first obtained for MDP with finite states and actions in (Williams \& Baird, 1990), and more recently in (Lan, 2023; Zhang et al., 2020). Finite time convergence for natural actor critic using a linear MDP assumption has been obtained in (Chen \& Zhao, 2022; Khodadadian et al., 2021; Xu et al., 2020b) with the best known sample complexity of $\tilde{\mathcal{O}}\left(\frac{1}{\epsilon^{2}(1-\gamma)^{4}}\right)$ (Xu et al., 2020a). Finite time sample complexity results are however, not available for Natural Actor Critic setups for general MDP where neural networks are used to represent the critic. (Fu et al., 2020) obtained asymptotic results for a variant of the Natural Actor Critic using a PPO update for the policy gradient step, but it forgoes the use of the 'clipped surrogate objective', which makes the algorithm unsuitable practically. The key related works here are summarized in Table 1.

## 3 Problem Formulation

We consider a discounted Markov Decision Process (MDP) given by the tuple $\mathcal{M}:=$ $(\mathcal{S}, \mathcal{A}, P, R, \gamma)$, where $\mathcal{S}$ is a bounded measurable state space, $\mathcal{A}$ is the finite set of actions. $P: \mathcal{S} \times \mathcal{A} \rightarrow \mathcal{P}(\mathcal{S})$ is the probability transition kernel $^{1}, R: \mathcal{S} \times \mathcal{A} \rightarrow \mathcal{P}\left(\left[0, R_{\max }\right]\right)$ is the reward kernel on the state action space with $R_{\text {max }}$ being the absolute value of the maximum reward, and $0<\gamma<1$ is the discount factor. A policy $\pi: \mathcal{S} \rightarrow \mathcal{P}(\mathcal{A})$ maps a state to a probability distribution over the action space. The action value function for a given policy $\pi$ is given by

$$
\begin{equation*}
Q^{\pi}(s, a)=\mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r\left(s_{t}, a_{t}\right) \mid s_{0}=s, a_{0}=a\right] \tag{1}
\end{equation*}
$$

where $r\left(s_{t}, a_{t}\right) \sim R\left(\cdot \mid s_{t}, a_{t}\right), a_{t} \sim \pi\left(\cdot \mid s_{t}\right)$ and $s_{t+1} \sim P\left(\cdot \mid s_{t}, a_{t}\right)$ for $t=\{0, \cdots, \infty\}$. For a discounted MDP, we define the optimal action value functions as

$$
\begin{equation*}
Q^{*}(s, a)=\sup _{\pi} Q^{\pi}(s, a), \quad \forall(s, a) \in \mathcal{S} \times \mathcal{A} \tag{2}
\end{equation*}
$$

A policy that achieves the optimal action-value functions is known as the optimal policy and is denoted as $\pi^{*}$. Similarly, we can define the value function as $V^{\pi}(s)=\mathbb{E}\left[\sum_{t=0}^{\infty} \gamma^{t} r^{\prime}\left(s_{t}, a_{t}\right) \mid s_{0}=s\right]$, and from the definition of $Q^{\pi}(s, a)$, it holds that $V^{\pi}(s)=\mathbb{E}_{a \sim \pi}\left[Q^{\pi}(s, a)\right]$. Similarly, we can define the optimal value function as $V^{*}(s)=\sup _{\pi} V^{\pi}(s), \quad \forall s \in \mathcal{S}$.

We define $\rho_{\nu}^{\pi}(s)$ as the stationary state distribution induced by the policy $\pi$ starting at state distribution $\nu$ and $\zeta_{\nu}^{\pi}(s, a)$ is the corresponding stationary state action distribution defined as $\zeta_{\nu}^{\pi}(s, a)=\rho_{\nu}^{\pi}(s) \pi(a \mid s)$. We further define $V^{\pi}(\nu)=\mathbb{E}_{s_{0} \sim \nu}\left[V^{\pi}\left(s_{0}\right)\right]$, where $\nu$ is an initial state distribution. We can define the visitation distribution as $d_{s_{0}}^{\pi}(s)=(1-\gamma) \sum_{t=0}^{\infty} \gamma^{t} \operatorname{Pr}^{\pi}\left(s_{t}=s \mid s_{o}\right)$. Here $\operatorname{Pr}^{\pi}\left(s_{t}=s \mid s_{o}\right)$ denotes the probability the state at time $t$ is $s$ given a starting state of $s_{o}$. Hence, we can write $d_{\rho}^{\pi}(s)=\mathbb{E}_{s_{o} \sim \rho}\left[d_{s_{0}}^{\pi}(s)\right]$. Finally for any measurable function $f: \mathcal{S} \times \mathcal{A} \rightarrow \mathbb{R}$ and a measure $\nu$ defined on $\mathcal{S} \times \mathcal{A}$ we define $\mathbb{E}(f)_{\nu}=\int_{\mathcal{S} \times \mathcal{A}} f d \nu$, ,
We additionally define the bellman operator for a policy $\pi$ on a function $Q: \mathcal{S} \times \mathcal{A} \rightarrow \mathbb{R}$ is defined as

$$
\begin{equation*}
\left(T^{\pi} Q\right)(s, a)=\mathbb{E}(r(s, a))+\gamma \int Q\left(s^{\prime}, \pi\left(s^{\prime}\right)\right) P\left(d s^{\prime} \mid s, a\right) \tag{3}
\end{equation*}
$$

Further, operator $P^{\pi}$ is defined as

$$
\begin{equation*}
P^{\pi} Q(s, a)=\mathbb{E}\left[Q\left(s^{\prime}, a^{\prime}\right) \mid s^{\prime} \sim P(\cdot \mid s, a), a^{\prime} \sim \pi\left(\cdot \mid s^{\prime}\right)\right] \tag{4}
\end{equation*}
$$

This is the one step Markov transition operator for policy $\pi$ for the Markov chain defined on $\mathcal{S} \times \mathcal{A}$ with the transition dynamics given by $S_{t+1} \sim P\left(\cdot \mid S_{t}, A_{t}\right)$ and $A_{t+1} \sim \pi\left(\cdot \mid S_{t+1}\right)$. It defines a distribution on the state action space after one transition from the initial state. Similarly, $P^{\pi_{1}} P^{\pi_{2}} \cdots P^{\pi_{m}}$ is the $m$-step Markov transition operator following policy $\pi_{t}$ at steps $1 \leq t \leq m$.

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## 4 Natural Actor Critic Algorithm Overview

We now describe our natural actor-critic (NAC) algorithm. In a natural policy gradient algorithm (Kakade, 2001a), the policy is parameterized as $\left\{\pi_{\lambda}, \lambda \in \Lambda\right\}$ and $\Lambda \subset \mathbb{R}^{d}$ where $d$ is a positive integer. We have $K$ total iterations of the Algorithm. At iteration $k$, the policy parameters are updated using a natural policy gradient step given by

$$
\begin{equation*}
\lambda_{k+1}=\lambda_{k}+\eta F_{\nu}^{\dagger}(\lambda) \nabla_{\lambda} V^{\pi_{\lambda}}(\nu) \tag{5}
\end{equation*}
$$

From the policy gradient theorem in (Sutton et al., 1999) we have

$$
\begin{align*}
\nabla_{\lambda_{k}} V^{\pi_{\lambda_{k}}}(\nu) & =\mathbb{E}_{s, a}\left(\nabla \log \left(\pi_{\lambda_{k}}\right)(a \mid s) Q^{\pi_{\lambda_{k}}}(s, a)\right)  \tag{6}\\
F_{\nu}\left(\lambda_{k}\right) & =\mathbb{E}_{s, a}\left[\nabla \log \pi_{\lambda_{k}}(a \mid s)\left(\nabla_{t} \log \pi_{\lambda_{k}}(a \mid s)\right)^{\top}\right] \tag{7}
\end{align*}
$$

where $s \sim d_{\nu}^{\pi_{\lambda_{k}}}, a \sim \pi_{\lambda_{k}}(. \mid s)$. From Sutton et al. (1999), the principle of compatible function approximation implies that we have

$$
\begin{align*}
F_{\nu}^{\dagger}\left(\lambda_{k}\right) \nabla_{\lambda_{k}} V^{\pi_{\lambda_{k}}}(\nu) & =\frac{1}{1-\gamma} w_{k}^{*}  \tag{8}\\
w_{k}^{*} & =\underset{w}{\arg \min } \mathbb{E}_{s, a}\left(A^{\pi_{\lambda_{k}}}(s, a)-w \nabla_{\lambda} \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right)^{2} \tag{9}
\end{align*}
$$

and $s \sim d_{\nu}^{\pi_{\lambda_{k}}}, a \sim \pi_{\lambda_{k}}(. \mid s)$ Here $\left(A^{\pi_{\lambda_{k}}}(s, a)=Q^{\pi_{\lambda_{k}}}(s, a)-V^{\pi_{\lambda_{k}}}(s)\right)$ and where $F^{\dagger}$ denotes the Moore-Penrose pseudo-inverse of the matrix $F$. For natural policy gradient algorithms such as in Agarwal et al. (2021) and Liu et al. (2020b) an estimate of $Q^{\pi_{\lambda_{k}}}$ (and from that an estimate of $A^{\pi_{\lambda_{k}}}(s, a)$ ) is obtained through a sampling procedure that requires on average $\left(\frac{1}{1-\gamma}\right)$ for each sample of $Q^{\pi_{\lambda_{k}}}$ (and thus $A^{\pi_{\lambda_{k}}}$ ). For the natural actor-critic setup, we maintain a parameterized estimate of the $Q$-function, which is updated at each step and is used to approximate $Q^{\pi_{\lambda_{k}}}$. In our case, a neural network with $L$ layers and at least $m$ neurons per layer is used to represent the $Q$ function, at each iteration $k$ of the algorithm, an estimate of its parameters is obtained by solving an optimization of the form

$$
\begin{equation*}
\underset{\theta \in \Theta}{\arg \min } \mathbb{E}_{s, a}\left(Q^{\pi_{\lambda_{k}}}-Q_{\theta}\right)^{2}, \tag{10}
\end{equation*}
$$

Where $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}, \Theta$ is the space of parameters for the neural networks and $Q_{\theta}$ is the neural network corresponding to the parameter $\theta$. This step is known as the critic step. A DQN like algorithm to get an estimate of $Q^{\pi_{\lambda_{k}}}$, as is done in practical implementations of the Natural Actor Critic like Wang et al. (2021a). We summarize the Natural Actor-Critic approach in Algorithm 1. It has one main for loop indexed by the iteration counter $k$. The first inner for loop indexed by $j$ is the loop where the critic step is performed. At a fixed iteration $k$ of the main for loop and iteration $j$ of the first inner for loop, we solve the following optimization problem

$$
\begin{equation*}
\underset{\theta \in \Theta}{\arg \min } \mathbb{E}_{s, a}\left(T^{\pi_{\lambda_{k}}} Q_{k, j-1}(s, a)-Q_{\theta}(s, a)\right)^{2} \tag{11}
\end{equation*}
$$

This is equivalent to the target network feature of the Deep Q Network(DQN) algorithm. For the inner loop at iteration $j$, the target is fixed to be $T^{\pi_{\lambda_{k}}} Q_{k, j-1}(s, a)$. The first inner for loop has a nested inner for loop indexed by $i$ where the optimization step for the current target is performed. The target network is updated at the end of the first inner loop. We note that the target network technique is applied in most real-world applications of natural actor critic with neural network critic as in (Wei et al., 2019). The first inner loop controls how many times the target network is updated. To get rid of the Markov dependence between the samples, the replay buffer technique is used wherein we randomly sample from the collected data instead of using it sequentially. For the sake of generality, we have not used this as our analysis will account for the Markov dependence between the samples.
The estimate of $w_{k}^{*}$ is obtained in the second inner for loop of Algorithm equation 1 indexed by $i$ where a gradient descent is performed for the loss function of the form given in equation 9 using the state action pairs sampled in the first inner for loop. Note that we do not have access to the true advantage function required for the critic update. Thus, we use the estimate of the $Q$ function obtained at the end of the first inner for loop to calculate the advantage function. After obtaining our estimate of the minimizer of equation 9 , we update the policy parameter using the stochastic gradient update step. Here, the state action pairs used are the same we sampled in the first inner for loop.

```
Algorithm 1 Natural Actor Critic with Neural Parametrization
Input: \(\mathcal{S}, \mathcal{A}, \gamma\), Time Horizon \(\mathrm{K} \in \mathcal{Z}\), Updates per time step \(\mathrm{J} \in \mathcal{Z}\),starting state sampling
distribution \(\nu\), Actor step sizes \(\beta_{i, k}, \forall k \in\{1, \cdots, K\}, i \in\{1, \cdots, n . J\}\), Critic step size \(\alpha\), policy
gradient step size \(\eta\),
    Initialize: \(\lambda_{0}=\{0\}^{d}\),
    for \(k \in\{1, \cdots, K\}\) do
        Initialize \(X=\varnothing, Q_{k}(s, a)=0 \quad \forall(s, a) \in \mathcal{S} \times \mathcal{A}\)
        for \(j \in\{1, \cdots, J\}\) do
            Sample \(s_{1}\) from \(\nu\) and \(a_{1}\) by following \(\pi_{\lambda_{k}}\)
            Initialize \(\theta_{0}\) using a standard Gaussian.
            for \(i \in\{1, \cdots, n\}\) do
                Sample the tuple \(s_{i+1}, a_{i+1}\) by following the policy \(\pi_{\lambda_{k}}\)
                Set \(y_{i}=r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k}\left(s_{i+1}, a_{i+1}\right)\),
                \(\theta_{i}=\theta_{i-1}+\alpha\left(y_{i}-Q_{\theta_{i}}\left(s_{i}, a_{i}\right)\right) \nabla Q_{\theta_{i}}\left(s_{i}, a_{i}\right)\)
            end for
            \(Q_{k}=Q_{\theta_{n}}\)
            Append the \(n\left(s_{i}, a_{i}\right)\) pairs to the data-set \(X\)
        end for
        Initialize \(w_{0}=0^{d}\)
        for \(i \in\{1, \cdots,|X|\}\) do
            \(A_{k}\left(s_{i}, a_{i}\right)=Q_{k}\left(s_{i}, a_{i}\right)-\sum_{a \in \mathcal{A}} \pi_{\lambda_{k}}\left(a \mid s_{i}\right) Q_{k}\left(s_{i}, a\right)\)
            \(w_{i}=w_{i-1}-\beta_{i, k}\left(w_{i} \cdot \nabla_{\lambda} \log \pi_{\lambda_{k}}\left(a_{i} \mid s_{i}\right)-A_{k}\left(s_{i}, a_{i}\right)\right) \nabla_{\lambda} \log \pi_{\lambda_{k}}\left(a_{i} \mid s_{i}\right)\)
        end for
        Update \(\lambda_{k+1}=\lambda_{k}+\eta w_{|X|}\)
    end for
    Output: \(\pi_{\lambda_{K+1}}\)
```


## 5 Global Convergence Result

### 5.1 ASSUMPTIONS

Before stating the main result, we formally describe the required assumptions in this subsection.
Assumption 1. For any $\lambda_{1}, \lambda_{2} \in \Lambda$ and $(s, a) \in(\mathcal{S} \times \mathcal{A})$ we have

$$
\begin{equation*}
\left\|\nabla \log \left(\pi_{\lambda_{1}}\right)(a \mid s)-\nabla \log \left(\pi_{\lambda_{2}}\right)(a \mid s)\right\|_{2} \leq \beta\left\|\lambda_{1}-\lambda_{2}\right\|_{2} \tag{12}
\end{equation*}
$$

where $\beta>0$.
Such assumptions have been utilized in prior policy gradient based works such as Agarwal et al. (2021); Liu et al. (2020b) and finite time analysis of NAC using linear critic such as Xu et al. (2020a). This assumption is satisfied for the softmax policy parameterization

$$
\begin{equation*}
\pi_{\lambda}(a \mid s)=\frac{\exp \left(f_{\lambda}(s, a)\right)}{\sum_{a^{\prime} \in \mathcal{A}} \exp \left(f_{\lambda}\left(s, a^{\prime}\right)\right)} \tag{13}
\end{equation*}
$$

where $f_{\lambda}(s, a)$ is a neural network with a smooth activation function. This is the most common form of the policy used in practice (Wei et al., 2019; Wang et al., 2021a). This assumption is also satisfied for Gaussian (Doya, 2000) and Boltzmann policies (Konda \& Tsitsiklis, 1999). Thus our analysis is more general than Fu et al. (2020) which is restricted to energy-based policies.
Assumption 2. For any $\lambda \in \Lambda$, let $\pi_{\lambda}$ be the corresponding policy, $\nu$ be the starting distribution over the state space, and let $\zeta_{\nu}^{\pi_{\lambda}}$ be the corresponding stationary state action distribution. We assume that there exists a positive integer $p$ such that for every positive integer $\tau$

$$
\begin{equation*}
d_{T V}\left(\mathbb{P}\left(\left(s_{\tau}, a_{\tau}\right) \in \cdot \mid\left(s_{0}, a_{0}\right)=(s, a)\right), \zeta_{\nu}^{\pi_{\lambda}}(\cdot)\right) \leq p \rho^{\tau}, \forall(s, a) \in \mathcal{S} \times \mathcal{A} \tag{14}
\end{equation*}
$$

This assumption implies that the Markov chain is geometrically mixing. Such assumption is widely used both in the analysis of stochastic gradient descent literature such as Doan (2022); Sun et al.
(2018), as well as finite time analysis of RL algorithms such as Xu et al. (2020a). In Fu et al. (2020), it is assumed that data can be sampled from the stationary distribution of a given policy. We note that this is not possible in practice. Instead, we can only sample from a Markov chain which has a stationary distribution as the desired distribution to sample from.
Assumption 3. For any fixed $\lambda \in \Lambda$ and $\theta \in \Theta$ we have

$$
\begin{equation*}
\min _{w} \mathbb{E}_{s, a \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}}\left(A_{\theta}(s, a)-w^{\top} \nabla \log \left(\pi_{\lambda}\right)(a \mid s)\right)^{2} \leq \epsilon_{\text {bias }} \tag{15}
\end{equation*}
$$

Similar assumptions are made in Fu et al. (2020), where this error is assumed to be zero when the critic has a linear function parameterization. In policy gradient works such as Liu et al. (2020b), the assumption replaces the parameterised estimate of the advantage function $A_{\theta}$ (which is known to us) with the true advantage function for policy $\pi_{\lambda}$ denoted by $A^{\pi_{\lambda}}$ (which is unknown to us). Doing so ignores the error that is incurred due to a mismatch in the actor and critic parameterization which is a critical aspect of a successful implementation of natural actor-critic algorithms. In Xu et al. (2020a), this assumption is implicit as this term is defined as a constant denoted by $\zeta_{\text {approx }}^{\text {actor }}$.
Assumption 4. For any fixed $\theta \in \Theta$ and $\lambda \in \Lambda$ we have

$$
\begin{equation*}
\min _{\theta_{1} \in \Theta} \mathbb{E}_{s, a \sim \zeta_{\nu}^{\pi_{\lambda}}}\left(Q_{\theta_{1}}(s, a)-T^{\pi_{\lambda}} Q_{\theta}(s, a)\right)^{2} \leq \epsilon_{\text {approx }} \tag{16}
\end{equation*}
$$

This assumption is key to the validity of the DQN step. Note that in works such as Xu et al. (2020a), an upper bound is placed on the approximation error when the function class (in that case linear functions) are used to approximate the unknown true value function (see term denoted as $\zeta_{\text {critic }}^{\text {approx }}$ ). Our assumption is weaker as we only require the class of neural network to be able to approximate the function obtained by applying the bellman operator to a neural network belonging to the same class.

### 5.2 Main Result

Theorem 1. Suppose Assumptions 1-4 hold and we have $\alpha=\Theta\left(\frac{1}{\operatorname{poly}(n, L) \cdot m}\right), \beta_{i, k}=\frac{2}{\mu_{k}(i+1)}$ where $\mu_{k}$ is the strong convexity parameter of the loss function in equation $9, \eta=\frac{1}{\sqrt{K}}$ and $m \geq$ $\mathcal{O}\left(K . J . \delta^{-1}\right)$ then from Algorithm 1 we obtain with probability at least $1-\delta$

$$
\begin{align*}
\min _{k \leq K}\left(V^{*}(\nu)-V^{\pi_{\lambda_{k}}}(\nu)\right) \leq \mathcal{O} & \left(\frac{1}{\sqrt{K}(1-\gamma)}\right)+\frac{1}{K(1-\gamma)} \sum_{k=1}^{K}\left(\mathcal{O}\left(\frac{\log (J \cdot n)}{J \cdot n}\right)+\mathcal{O}\left(\gamma^{J}\right)\right) \\
& +\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} \sum_{j=0}^{J-1}\left(\mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{n}+\mathcal{O}\left(\frac{1}{\sqrt{n}}\right)\right) \\
& +\frac{1}{1-\gamma}\left(\mathcal{O}\left(\epsilon_{\text {bias }}\right)+\mathcal{O}\left(\sqrt{\epsilon_{\text {approx }}}\right)\right) \tag{17}
\end{align*}
$$

Hence, for $K=\mathcal{O}\left(\epsilon^{-2}(1-\gamma)^{-2}\right), J=\mathcal{O}\left(\log \left(\frac{1}{\epsilon}\right)\right), n=\tilde{\mathcal{O}}\left(\epsilon^{-2}(1-\gamma)^{-2}\right), m \geq \mathcal{O}\left(\epsilon^{-2} . \delta^{-1}\right)$

$$
\begin{equation*}
\min _{k \leq K}\left(V^{*}(\nu)-V^{\pi_{\lambda_{k}}}(\nu)\right) \leq \epsilon+\frac{1}{1-\gamma}\left(\epsilon_{\text {bias }}+\left(\sqrt{\epsilon_{\text {approx }}}\right)\right) \tag{18}
\end{equation*}
$$

which implies a sample complexity of $K \cdot J \cdot n=\tilde{\mathcal{O}}\left(\epsilon^{-4}(1-\gamma)^{-4}\right)$.
Remark 1: We note that there are seven terms on the right-hand side of equation 17. The first term is a consequence of the smoothness property of the actor parameterization. The second term is the error incurred in estimating $w_{k}^{*}$. The third term is the error incurred due to the inherent randomness of the system during each critic update step, in Farahmand et al. (2010) this was known as the statistical error. The fourth term on the right is the error incurred in fitting the data at each fixed target in the critic step. The fifth term on the right is the error incurred due to a lack of knowledge of the transition matrix. The sixth term $\epsilon_{b i a s}$ represents the minimum possible attainable value of the loss function in the actor step. It is also a measure of how compatible are the architecture of the actor and critic. In Wang et al. (2019) it is shown that for an over-parameterized neural network used to represent both
the actor and critic this error is zero. The term $\epsilon_{\text {approx }}$ is a measure of how well the class of neural networks we use to represent the critic can approximate a function obtained by applying the bellman operator to a function from that same class. Works such as Fan et al. (2020); Chen \& Jiang (2019) set this error to zero. The requirement on the minimum number of neurons $m$ in each layer of the critic network can be be thought of as a consequence of the universal approximation property which states that sufficiently wide neural networks(even those with a single hidden layer) can approximate any continuous function with arbitrary accuracy.
Remark 2: Our sample complexity when compared to the existing state of the art sample complexity bound for natural policy gradient with non-linear policy parameterization of $\tilde{\mathcal{O}}\left(\epsilon^{-3}(1-\gamma)^{-6}\right)$ achieved in Liu et al. (2020a) reveals a key insight. Note that our bound is worse off in terms of $\epsilon$ by a factor of $\epsilon^{-1}$. This is due to the fact that we have to obtain an estimate of the critic parameters while the natural policy gradient does not. We can see this in our result from the fifth term on the right hand side of equation 17 which is $\mathcal{O}\left(n^{-\frac{1}{2}}\right)$ which is from the critic optimization step. The natural policy gradient algorithm requires on average $(1-\gamma)^{-1}$ state action samples for every sample of $Q(s, a)$. This is reflected in our results as our error bounds are better in terms of $(1-\gamma)$ by a factor of $(1-\gamma)^{-2}$. We discuss this detail in Appendix E.

Remark 3: Note the presence of the probability term for our convergence result. This term is present due to the fact that the optimization for the critic step is non-convex, hence convergence can only be guaranteed with a high probability. We show in the Appendix F that if the critic is represented by a two layer neural network with ReLU activation, using the convex reformulation as laid out in Mishkin et al. (2022), a deterministic upper bound on the error can be obtained.

## 6 Proof Sketch of Theorem 1

The proof is split into two stages. In the first stage, we demonstrate how the difference in value functions is upper bounded as a function of the errors incurred till the final step $K$. The second part is to upper bound the different error components.
Upper Bounding Error in Separate Error Components: We use the smoothness property assumed in Assumption 1 to obtain a bound on the expectation of the difference between our estimated value function and the optimal value function.

$$
\begin{equation*}
\min _{k \in\{1, \cdots, K\}} V^{*}(\nu)-V^{\pi_{\lambda_{K}}}(\nu) \leq \frac{\log (|\mathcal{A}|)}{K \eta(1-\gamma)}+\frac{\eta \beta W^{2}}{2(1-\gamma)}+\frac{1}{K} \sum_{k=1}^{K} \frac{e r r_{k}}{1-\gamma}, \tag{19}
\end{equation*}
$$

where

$$
\begin{equation*}
\operatorname{err}_{k}=\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}(. \mid s)}\left(\left|A^{\pi_{\lambda_{k}}}-w^{k}(s, a) \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right|\right) \tag{20}
\end{equation*}
$$

where $W$ is a constant such that $\left\|w^{k}\right\|_{2} \leq W \forall k$, where $k$ denotes the iteration of the outer for loop of Algorithm 1. We split the term in equation 20 into the errors incurred due to the actor and critic step as follows

$$
\begin{align*}
\operatorname{err}_{k} & =\mathbb{E}_{s, a}\left(\left|A^{\pi_{\lambda_{k}}}-w^{k} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right|\right)  \tag{21}\\
& \leq \underbrace{\mathbb{E}_{s, a}\left(\left|A^{\pi_{\lambda_{k}}}-A_{k, J}\right|\right)}_{I}+\underbrace{\mathbb{E}_{s, a}\left(\left|A_{k, J}-w^{k} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right|\right)}_{I I} . \tag{22}
\end{align*}
$$

Note that $I$ is the difference between the true $A^{\pi_{\lambda_{k}}}$ function corresponding to the policy $\pi_{\lambda_{k}}$ and $A_{k, J}$ is our estimate. This estimation is carried out in the first inner for loop of Algorithm 1. Thus $I$ is the error incurred in the critic step. $I I$ is the error incurred in the estimation of the actor update. This is incurred in the stochastic gradient descent steps in the second inner for loop of Algorithm 1.

Upper Bounding Error in Critic Step: For each iteration $k$ of the Algorithm 1. We show that minimizing $I$ is equivalent to solving the following problem

$$
\begin{equation*}
\underset{\theta \in \Theta}{\arg \min } \mathbb{E}_{s, a}\left(Q^{\pi_{\lambda_{k}}}-Q_{\theta}\right)^{2} \tag{23}
\end{equation*}
$$

where $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}$. We recreate the result for the value function from Lemmas 2 of Munos (2003) for the action value function $Q$ to obtain

$$
\begin{equation*}
\mathbb{E}_{s, a}\left|Q^{\pi_{\lambda_{k}}}-Q_{k, J}\right| \leq \sum_{j=1}^{J-1} \gamma^{J-j-1}\left(P^{\pi_{\lambda_{k}}}\right)^{J-j-1} \mathbb{E}\left|\epsilon_{k, j}\right|+\gamma^{J}\left(\frac{R_{\max }}{1-\gamma}\right) \tag{24}
\end{equation*}
$$

where $\epsilon=T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}$ is the Bellman error incurred at iteration $j$ of the first inner for loop and iteration $k$ of the outer for loop of Algorithm 1. The first term on the right hand side is called as the algorithmic error, which depends on how good our approximation of the Bellman error is. The second term on the right hand side is called as the statistical error, which is the error incurred due to the random nature of the system. Intuitively, the Bellman error depends on how much data is collected at each iteration, how efficient our solution to the optimization step is to the true solution, and how well our function class can approximate $T^{\pi_{\lambda_{k}}} Q_{k, j-1}$. Building upon this intuition, we split $\epsilon$ into four different components as follows.

$$
\begin{align*}
\epsilon_{k, j} & =T^{\pi_{\lambda_{k}} Q_{k, j-1}-Q_{k, j}} \\
& =\underbrace{T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}}_{\epsilon_{k, j}^{1}}+\underbrace{Q_{k, j}^{1}-Q_{k, j}^{2}}_{\epsilon_{k, j}^{2}}+\underbrace{Q_{k, j}^{2}-Q_{k, j}^{3}}_{\epsilon_{k, j}^{3}}+\underbrace{Q_{k, j}^{3}-Q_{k, j}}_{\epsilon_{k, j}^{4}} \\
& =\epsilon_{k, j}^{1}+\epsilon_{k, j}^{2}+\epsilon_{k, j}^{3}+\epsilon_{k, j}^{4}, \tag{25}
\end{align*}
$$

We now define the terms introduced above. We first define the various $Q$-functions which we can approximate in decreasing order of the accuracy and then define the corresponding errors.

We start by defining the best possible approximation of the function $T^{\pi_{\lambda_{k}}} Q_{k, j-1}$ possible from the class of neural networks with smooth activation functions, with respect to the expected square from the true ground truth $T^{\pi_{\lambda_{k}}} Q_{k, j-1}$.
Definition 1. For iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, we define

$$
\begin{equation*}
Q_{k, j}^{1}=\underset{Q_{\theta}, \theta \in \Theta}{\arg \min } \mathbb{E}\left(Q_{\theta}(s, a)-T^{\pi_{\lambda_{k}}} Q_{k, j-1}(s, a)\right)^{2} \tag{26}
\end{equation*}
$$

where $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)$.
Note that we do not have access to the transition probability kernel $P$, hence we do not know $T^{\pi \lambda_{k}}$. To alleviate this, we use the observed next state and actions instead. Using this, we define $Q_{k, j}^{2}$ as,
Definition 2. For iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, we define

$$
\begin{equation*}
Q_{k, j}^{2}=\underset{Q_{\theta}, \theta \in \Theta}{\arg \min } \mathbb{E}\left(Q_{\theta}(s, a)-\left(r^{\prime}(s, a)+\gamma Q_{k, j-1}\left(s^{\prime}, a^{\prime}\right)\right)^{2},\right. \tag{27}
\end{equation*}
$$

where $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a), s^{\prime} \sim P\left(s^{\prime} \mid s, a\right), r^{\prime}(\cdot \mid s, a) \sim R(\cdot \mid s, a)$ and $a^{\prime} \sim \pi^{\lambda_{k}}\left(. \mid s^{\prime}\right)$
To obtain $Q_{k, j}^{2}$, we still need to compute the true expected value in Equation 27. However, we still do not know the transition function $P$. To remove this limitation, we use sampling. Consider the set of $n$ state-action pairs sampled by starting from a state action distribution $\nu$ and following policy $\pi^{\lambda_{k}}$, using which we define $Q_{k, j}^{3}$ as,
Definition 3. For the set of $n$ state action pairs sampled in iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1 we define

$$
\begin{equation*}
Q_{k, j}^{3}=\underset{Q_{\theta}, \theta \in \Theta}{\arg \min } \frac{1}{n} \sum_{i=1}^{n}\left(Q_{\theta}\left(s_{i}, a_{i}\right)-\left(r\left(s_{i}, a_{i}\right)+\gamma Q_{k, j-1}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2}, \tag{28}
\end{equation*}
$$

$Q_{k, j}^{3}$ is the best possible approximation for $Q$-value function which minimizes the sample average of the square loss functions with the target values as $\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j-1}\left(s_{i+1}, a_{i+1}\right)\right)$. In other words this is the optimal solution for fitting the observed data.

We now defined the errors using the $Q$ functions just defined. We start by defining the approximation error which represents the difference between the function $T^{\pi_{\lambda_{k}}} Q_{j-1}$ and its best approximation possible from the class of neural networks used for critic parametrization denoted by $Q_{k, j}^{1}$.

Definition 4 (Approximation Error). For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, we define, $\epsilon_{k, j}^{1}=T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}$, where $Q_{k, j-1}$ is the estimate of the $Q$ function at iteration $k$ of the outer for loop and iteration $j-1$ of the first inner for loop of Algorithm 1.

This error is a measure of the approximation power of the class of neural networks we use to represent the critic. We upper bound this error in lemma 3 in Appendix B.
We also define Estimation Error which denotes the error between the best approximation of $T^{\pi_{\lambda_{k}}} Q_{k, j-1}$ possible from the class of neural networks denoted by $Q_{k, j}^{1}$ and the minimizer of the loss function in equation 27 denoted $Q_{k, j}^{2}$.
Definition 5 (Estimation Error). For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, we define, $\epsilon_{k, j}^{2}=Q_{k, j}^{1}-Q_{k, j}^{2}$.

We demonstrate that this error is zero in lemma 4 in Appendix B.
We now define Sampling error which denotes the difference between the minimizer of expected loss function in equation 27 denoted by $Q_{k, j}^{2}$ and the minimizer of the empirical loss function in equation 28 denoted by $Q_{k, j}^{3}$. We can see that intuitively, the more samples we have the closer these two functions will be. We use Rademacher complexity results to upper bound this error.
Definition 6 (Sampling Error). For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, we define, $\epsilon_{k, j}^{3}=Q_{k, j}^{3}-Q_{k, j}^{2}$.

An upper bound on this error is established in 5 in Appendix B.
Lastly, we define optimization error which denotes the difference between the minimizer of the empirical square loss function, $Q_{k_{3}}$, and our estimate of this minimizer that is obtained from the gradient descent algorithm.
Definition 7 (Optimization Error). For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, we define, $\epsilon_{k}^{4}=Q_{k, j}^{3}-Q_{k, j}$. Here $Q_{k, j}$ is our estimate of the $Q$ function at iteration $k$ of Algorithm 1 and iteration $j$ of the first inner loop of Algorithm 1.

The upper bound on these error terms is established in lemma 6 in Appendix B.
Upper Bounding Error in Actor Step: Note that we require the minimization of the term $\mathbb{E}_{s, a}\left(A_{k, J}-w^{k} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right)$. Here the expectation is with respect to stationary state action distribution corresponding to $\pi_{\lambda_{k}}$. But we do not have samples of states action pairs from the stationary distribution with respect to the policy $\pi_{\lambda_{k}}$, we only have samples from the Markov chain induced by the policy $\pi_{\lambda_{k}}$. We thus refer to the theory in Doan (2022) and Assumption 3 to upper bound the error incurred.

For the error incurred in the actor update we define the related loss function as
Definition 8. For iteration $k$ of the outer for loop of Algorithm 1 , we define $w_{k}$ as the estimate of the minima of the loss function given by $\mathbb{E}_{(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}(s, a)}}\left(A_{k, J}(s, a)-(w) \nabla \log \left(\pi_{\lambda_{k}}\right)(a \mid s)\right)^{2}$ obtained at the end of the second inner for loop of Algorithm 1. We further define the true minima as

$$
\begin{equation*}
w_{k}^{*}=\underset{w}{\arg \min } \mathbb{E}_{(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)}\left(A_{k, J}(s, a)-(w) \nabla \log \left(\pi_{\lambda_{k}}\right)(a \mid s)\right)^{2} \tag{29}
\end{equation*}
$$

For finding the estimate $w_{k}$, we re-use the state action pairs sampled in the first inner for loop of Algorithm 1. The difference between our estimate $w_{k}$ and the $w_{k}^{*}$ (which is also the minimizer of $I I)$ is then used to upper bound the difference between the value of $I I$ at our estimate $w_{k}$ and the minimum possible value of $I I$ achieved at $w_{k}^{*}$ which is upper bounded using Assumption 3. Details of this are given in lemma 7 in Appendix B.

## 7 CONCLUSIONS

In this paper, we study a natural actor critic algorithm with a neural network used to represent both the actor and the critic and find the sample complexity guarantees for the algorithm. We show that our approach achieves a sample complexity of $\tilde{\mathcal{O}}\left(\epsilon^{-4}(1-\gamma)^{-4}\right)$. This demonstrates the first approach for achieving sample complexity beyond linear MDP assumptions for the critic.

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## Appendix

## A Supplementary lemmas

Here we provide some definitions and results that will be used to prove the lemmas stated in the paper.
Definition 9. For a given set $Z \subset \mathbb{R}^{n}$, we define the Rademacher complexity of the set $Z$ as

$$
\begin{equation*}
\operatorname{Rad}(Z)=\mathbb{E}\left(\sup _{z \in Z} \frac{1}{n} \sum_{i=1}^{d} \Omega_{i} z_{i}\right) \tag{30}
\end{equation*}
$$

where $\Omega_{i}$ is random variable such that $P\left(\Omega_{i}=1\right)=\frac{1}{2}, P\left(\Omega_{i}=-1\right)=\frac{1}{2}$ and $z_{i}$ are the coordinates of $z$ which is an element of the set $Z$
Lemma 1. Consider a set of observed data denoted by $z=\left\{z_{1}, z_{2}, \cdots z_{n}\right\} \in \mathbb{R}^{n}$, a parameter space $\Theta$, a loss function $\{l: \mathbb{R} \times \Theta \rightarrow \mathbb{R}\}$ where $0 \leq l(\theta, z) \leq 1 \forall(\theta, z) \in \Theta \times \mathbb{R}$. The empirical risk for a set of observed data as $R(\theta)=\frac{1}{n} \sum_{i=1}^{n} l\left(\theta, \overline{z_{i}}\right)$ and the population risk as $r(\theta)=\mathbb{E} l\left(\theta, \tilde{z_{i}}\right)$, where $\tilde{z_{i}}$ is a co-ordinate of $\tilde{z}$ sampled from some distribution over $Z$.
We define a set of functions denoted by $\mathcal{L}$ as

$$
\begin{equation*}
\mathcal{L}=\{z \in Z \rightarrow l(\theta, z) \in \mathbb{R}: \theta \in \Theta\} \tag{31}
\end{equation*}
$$

Given $z=\left\{z_{1}, z_{2}, z_{3} \cdots, z_{n}\right\}$ we further define a set $\mathcal{L} \circ z$ as

$$
\begin{equation*}
\mathcal{L} \circ z=\left\{\left(l\left(\theta, z_{1}\right), l\left(\theta, z_{2}\right), \cdots, l\left(\theta, z_{n}\right)\right) \in \mathbb{R}^{n}: \theta \in \Theta\right\} \tag{32}
\end{equation*}
$$

Then, we have

$$
\begin{equation*}
\mathbb{E} \sup _{\theta \in \Theta}|\{r(\theta)-R(\theta)\}| \leq 2 \mathbb{E}(\operatorname{Rad}(\mathcal{L} \circ z)) \tag{33}
\end{equation*}
$$

If the data is of the form $z_{i}=\left(x_{i}, y_{i}\right), x \in X, y \in Y$ and the loss function is of the form $l\left(a_{\theta}(x), y\right)$, is L lipschitz and $a_{\theta}: \Theta \times X \rightarrow \mathbb{R}$, then we have

$$
\begin{equation*}
\mathbb{E} \sup _{\theta \in \Theta}|\{r(\theta)-R(\theta)\}| \leq 2 L \mathbb{E}\left(\operatorname{Rad}\left(\mathcal{A} \circ\left\{x_{1}, x_{2}, x_{3}, \cdots, x_{n}\right\}\right)\right) \tag{34}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{A} \circ\left\{x_{1}, x_{2}, \cdots, x_{n}\right\}=\left\{\left(a\left(\theta, x_{1}\right), a\left(\theta, x_{2}\right), \cdots, a\left(\theta, x_{n}\right)\right) \in \mathbb{R}^{n}: \theta \in \Theta\right\} \tag{35}
\end{equation*}
$$

The detailed proof of the above statement is given in (Rebeschini, 2022) ${ }^{2}$. The upper bound for $\mathbb{E} \sup _{\theta \in \Theta}(\{r(\theta)-R(\theta)\})$ is proved in the aformentioned reference. However, without loss of generality the same proof holds for the upper bound for $\mathbb{E} \sup _{\theta \in \Theta}(\{R(\theta)-r(\theta)\})$. Hence the upper bound for $\mathbb{E} \sup _{\theta \in \Theta}|\{r(\theta)-R(\theta)\}|$ can be established.

Lemma 2. Consider three random random variable $x \in \mathcal{X}$ and $y, y \in \mathcal{Y}$. Let $\mathbb{E}_{x, y}, \mathbb{E}_{x}$ and $\mathbb{E}_{y \mid x}$, $\mathbb{E}_{y^{\prime} \mid x}$ denote the expectation with respect to the joint distribution of $(x, y)$, the marginal distribution of $x$, the conditional distribution of $y$ given $x$ and the conditional distribution of $y^{\prime}$ given $x$ respectively. Let $f_{\theta}(x)$ denote a bounded measurable function of $x$ parameterised by some parameter $\theta$ and $g(x, y)$ be bounded measurable function of both $x$ and $y$.
Then we have

$$
\begin{equation*}
\underset{f_{\theta}}{\arg \min } \mathbb{E}_{x, y}\left(f_{\theta}(x)-g(x, y)\right)^{2}=\underset{f_{\theta}}{\arg \min }\left(\mathbb{E}_{x}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)^{2}\right) \tag{36}
\end{equation*}
$$

[^1]Proof. Denote the left hand side of Equation equation 36 as $\mathbb{X}_{\theta}$, then add and subtract $\mathbb{E}_{y \mid x}(g(x, y) \mid x)$ to it to get

$$
\begin{align*}
\mathbb{X}_{\theta}= & \underset{f_{\theta}}{\arg \min }\left(\mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)+\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)-g(x, y)\right)^{2}\right)  \tag{37}\\
= & \underset{f_{\theta}}{\arg \min }\left(\mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)^{2}+\mathbb{E}_{x, y}\left(y-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)^{2}\right. \\
& \left.\quad-2 \mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)\left(g(x, y)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)\right) . \tag{38}
\end{align*}
$$

Consider the third term on the right hand side of Equation equation 38

$$
\begin{align*}
& 2 \mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)\left(g(x, y)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right) \\
& \quad= 2 \mathbb{E}_{x} \mathbb{E}_{y \mid x}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)\left(g(x, y)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)  \tag{39}\\
& \quad= 2 \mathbb{E}_{x}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right) \mathbb{E}_{y \mid x}\left(g(x, y)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)  \tag{40}\\
&= 2 \mathbb{E}_{x}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)\left(\mathbb{E}_{y \mid x}(g(x, y))-\mathbb{E}_{y \mid x}\left(\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)\right)  \tag{41}\\
& \quad= 2 \mathbb{E}_{x}\left(f_{\theta}(x)-\mathbb{E}(y \mid x)\right)\left(\mathbb{E}_{y \mid x}(g(x, y))-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)  \tag{42}\\
&= 0 \tag{43}
\end{align*}
$$

Equation equation 39 is obtained by writing $\mathbb{E}_{x, y}=\mathbb{E}_{x} \mathbb{E}_{y \mid x}$ from the law of total expectation. Equation equation 40 is obtained from equation 39 as the term $f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)$ is not a function of $y$. Equation equation 41 is obtained from equation 40 as $\mathbb{E}_{y \mid x}\left(\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)=$ $\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)$ because $\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)$ is not a function of $y$ hence is constant with respect to the expectation operator $\mathbb{E}_{y \mid x}$.
Thus plugging in value of $2 \mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)\left(g(x, y)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)$ in Equation equation 38 we get

$$
\begin{align*}
& \underset{f_{\theta}}{\arg \min } \mathbb{E}_{x, y}\left(f_{\theta}(x)-g(x, y)\right)^{2}=\underset{f_{\theta}}{\arg \min }\left(\mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}_{x, y^{\prime}}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)^{2}\right. \\
&\left.+\mathbb{E}_{x, y}\left(g(x, y)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)^{2}\right) \tag{44}
\end{align*}
$$

Note that the second term on the right hand side of Equation equation 44 des not depend on $f_{\theta}(x)$ therefore we can write Equation equation 44 as

$$
\begin{equation*}
\underset{f_{\theta}}{\arg \min } \mathbb{E}_{x, y}\left(f_{\theta}(x)-g(x, y)\right)^{2}=\underset{f_{\theta}}{\arg \min }\left(\mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)^{2}\right) \tag{45}
\end{equation*}
$$

Since the right hand side of Equation equation 45 is not a function of $y$ we can replace $\mathbb{E}_{x, y}$ with $\mathbb{E}_{x}$ to get

$$
\begin{equation*}
\underset{f_{\theta}}{\arg \min } \mathbb{E}_{x, y}\left(f_{\theta}(x)-g(x, y)\right)^{2}=\underset{f_{\theta}}{\arg \min }\left(\mathbb{E}_{x}\left(f_{\theta}(x)-\mathbb{E}_{y^{\prime} \mid x}\left(g\left(x, y^{\prime}\right) \mid x\right)\right)^{2}\right) \tag{46}
\end{equation*}
$$

## B Supporting Lemmas

We will now state the key lemmas that will be used for finding the sample complexity of the proposed algorithm.

Lemma 3. For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, the approximation error denoted by $\epsilon_{k, j}^{1}$ in Definition 4, we have

$$
\begin{equation*}
\mathbb{E}\left(\left|\epsilon_{k, j}^{1}\right|\right) \leq \sqrt{\epsilon_{\text {approx }}} \tag{47}
\end{equation*}
$$

Where the expectation is with respect to and $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)$
Proof Sketch: We use Assumption 4 and the definition of the variance of a random variable to obtain the required result. The detailed proof is given in Appendix D.1.
Lemma 4. For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, $Q_{k, j}^{1}=Q_{k, j}^{2}$, or equivalently $\epsilon_{k, j}^{2}=0$

Proof Sketch: We use Lemma 2 in Appendix A and use the definitions of $Q_{k, j}^{1}$ and $Q_{k, j}^{2}$ to prove this result. The detailed proof is given in Appendix D.2.
Lemma 5. For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, if the number of state action pairs sampled are denoted by $n_{k, j}$, then the error $\epsilon_{k, j}^{3}$ defined in Definition 6 is upper bounded as

$$
\begin{equation*}
\mathbb{E}\left(\left|\epsilon_{k, j}^{3}\right|\right) \leq \tilde{\mathcal{O}}\left(\frac{1}{\sqrt{n}}\right) \tag{48}
\end{equation*}
$$

Where the expectation is with respect to and $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)$
Proof Sketch: First we note that For a given iteration $k$ of Algorithm 1 and iteration $j$ of the first for loop of Algorithm 1, $\mathbb{E}\left(R_{X_{k, j}, Q_{k, j-1}}(\theta)\right)=L_{Q_{j, k-1}}(\theta)$ where $R_{X_{k, j}, Q_{j, k-1}}(\theta)$ and $L_{Q_{j, k-1}}(\theta)$ are defined in Appendix D.3. We use this to get a probabilistic bound on the expected value of $\left|\left(Q_{j, k}^{2}\right)-\left(Q_{j, k}^{3}\right)\right|$ using Rademacher complexity theory when the samples are drawn from an ergodic Markov chain. The detailed proof is given in Appendix D.3. Note the presence of the $\log \left(\log \left(n_{k, j}\right)\right)$ term is due to the fact that the state action samples belong to a Markov Chain.
Lemma 6. For a given iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop of Algorithm 1, let the number of steps of the gradient descent performed by Algorithm 1, denoted by $T_{k, j}$, the minimum number of neurons $m$ satisfy $m \geq \mathcal{O}\left(\delta^{-1}\right)$ and the gradient descent step size $\alpha$ satisfy

$$
\begin{equation*}
\alpha=\Theta\left(\frac{1}{\operatorname{poly}(n, L) \cdot m}\right), \tag{49}
\end{equation*}
$$

Then with probability at least $1-\delta$ the error $\epsilon_{k_{4}}$ defined in Definition 7 is upper bounded as

$$
\begin{equation*}
\mathbb{E}\left(\left|\epsilon_{k, j}^{4}\right|\right) \leq \mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{T_{k, j}} \tag{50}
\end{equation*}
$$

Where the expectation is with respect to $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)$.
Proof Sketch: We use This is Theorem 5 from Zhu \& Xu (2021) to prove this lemma. The detailed proof is given in Appendix D.4.
Lemma 7. For a given iteration $k$ of the outer for loop of Algorithm 1, if the number of samples of the state action pairs sampled at each iteration of the first inner for loop are denoted by $n$ and $\beta_{i}$ be the step size in the gradient descent at iteration $i$ of the second inner for loop of Algorithm 1 which satisfies

$$
\begin{equation*}
\beta_{i}=\frac{2}{\mu_{k}(i+1)} \tag{51}
\end{equation*}
$$

where $\mu_{k}$ is the strong convexity parameter of the loss function $F_{k}=\mathbb{E}_{s, a \sim \zeta_{\nu}^{\pi_{\lambda_{k}}(s, a)}}\left(A_{k, J}-\right.$ $\left.w^{k} \nabla \log \left(\pi^{\lambda_{k}}(a \mid s)\right)\right)^{2}$. Then, after J.n iterations of gradient descent it holds that,

$$
\begin{equation*}
\left(F_{k}\left(w_{i}\right)\right) \leq \tilde{\mathcal{O}}\left(\frac{\log (J . n)}{J . n}\right)+F_{k}^{*} \tag{52}
\end{equation*}
$$

where $F_{k}^{*}=\arg \min _{w} F_{k}(w)$.

Proof Sketch: Note that we don not have access to state action samples belonging to the stationary state action distribution corresponding to the policy $\pi_{\lambda_{k}}$. We only have access to samples from Markov chain with the same stationary state action distribution. To account for this, we use the results in Doan (2022) and obtain the difference between the optimal loss function and the loss function obtained by performing stochastic gradient descent with samples from a Markov chain. The detailed proof is given in Appendix

## C Proof of Theorem 1

Proof. From Assumption 1, we have

$$
\begin{align*}
\log \frac{\pi_{\lambda_{k+1}}(a \mid s)}{\pi_{\lambda_{k}}(a \mid s)} & \geq \nabla_{\lambda_{k}} \log \pi_{\lambda_{k}}(a \mid s) \cdot\left(\lambda^{k+1}-\lambda^{k}\right)-\frac{\beta}{2}\left\|\lambda^{k+1}-\lambda^{k}\right\|_{2}^{2}  \tag{53}\\
& =\eta \log \pi_{\lambda_{k}}(a \mid s) \cdot w^{k}-\eta^{2} \frac{\beta}{2}\left\|w^{k}\right\|_{2}^{2} \tag{54}
\end{align*}
$$

From the definition of KL divergence and from the performance difference lemma from (Kakade \& Langford, 2002) we have

$$
\begin{align*}
\left.\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}}\left(K L\left(\pi^{*} \| \pi^{\lambda_{k}}\right)-\pi^{*}| | \pi^{\lambda_{k+1}}\right)\right)= & \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}} \mathbb{E}_{a \sim \pi^{*}(. \mid s)}\left[\log \frac{\pi^{\lambda_{k+1}}(a \mid s)}{\pi_{\lambda_{k}}(a \mid s)}\right]  \tag{55}\\
\geq & \eta \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}} \mathbb{E}_{a \sim \pi^{*}(. \mid s)}\left[\nabla_{\lambda_{k}} \log \pi_{\lambda_{k}}(a \mid s) \cdot w^{k}\right]-\eta^{2} \frac{\beta}{2}\left\|w^{k}\right\|_{2}^{2}  \tag{56}\\
= & \eta \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}} \mathbb{E}_{a \sim \pi^{*}(. \mid s)}\left[A^{\pi^{\lambda_{k}}}\right]-\eta^{2} \frac{\beta}{2}\left\|w^{k}\right\|_{2}^{2} \\
& +\eta \mathbb{E}_{s \sim d_{\nu}^{* *}} \mathbb{E}_{a \sim \pi^{*}(. \mid s)}\left[\nabla_{\lambda_{k}} \log \pi_{\lambda_{k}}(a \mid s) \cdot w^{k}-A^{\pi^{\lambda_{k}}}(s, a)\right]  \tag{57}\\
= & (1-\gamma) \eta\left(V^{\pi^{*}}(\nu)-V^{k}(\nu)\right)-\eta^{2} \frac{\beta}{2}\left\|w^{k}\right\|_{2}^{2}-\eta \cdot \operatorname{err}_{k} \tag{58}
\end{align*}
$$

Equation equation 58 is obtained from Equation equation 58 using the performance difference lemma form (Kakade \& Langford, 2002) where $A^{\pi^{\lambda k}}$ is the advantage function to the corresponding to the policy $\pi^{\lambda_{k}}$.

Rearranging, we get

$$
\begin{equation*}
\left(V^{\pi^{*}}(\nu)-V^{k}(\nu)\right) \leq \frac{1}{1-\gamma}\left(\frac{1}{\eta} \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}}\left(K L\left(\pi^{*} \| \pi^{\lambda_{k}}\right)-K L\left(\pi^{*} \| \pi^{\lambda_{k+1}}\right)\right)+\eta^{2} \frac{\beta}{2} \cdot W^{2}+\eta \cdot \operatorname{err}_{k}\right) \tag{59}
\end{equation*}
$$

Summing from 1 to $K$ and dividing by $K$ we get

$$
\begin{align*}
\frac{1}{K} \sum_{k=1}^{K}\left(V^{\pi^{*}}(\nu)-V^{k}(\nu)\right) & \leq\left(\frac{1}{1-\gamma}\right) \frac{1}{K} \sum_{k=1}^{K}\left(\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}}\left(K L\left(\pi^{*} \| \pi^{\lambda_{k}}\right)-K L\left(\pi^{*} \| \pi^{\lambda_{k+1}}\right)\right)+\eta \cdot \operatorname{err}_{k}\right) \\
& +\left(\frac{1}{1-\gamma}\right) \eta^{2} \frac{\beta}{2} \cdot W^{2}  \tag{60}\\
& \leq \frac{1}{\eta(1-\gamma)} \frac{1}{K} \mathbb{E}_{s \sim \tilde{d}}\left(K L\left(\pi^{*} \| \pi^{\lambda_{0}}\right)\right)+\frac{\eta \beta \cdot W^{2}}{2(1-\gamma)}+\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} e r r_{k}  \tag{61}\\
& \leq \frac{\log (|\mathcal{A}|)}{K \eta(1-\gamma)}+\frac{\eta^{2} \beta \cdot W^{2}}{2(1-\gamma)}+\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} \operatorname{err}_{k} \tag{62}
\end{align*}
$$

If we set $\eta=\frac{1}{\sqrt{K}}$ in Equation equation 62 we get

$$
\begin{equation*}
\frac{1}{K} \sum_{k=1}^{K}\left(V^{\pi^{*}}(\nu)-V^{k}(\nu)\right) \leq \frac{1}{\sqrt{K}}\left(\frac{2 \log (|\mathcal{A}|)+\beta \cdot W^{2}}{2(1-\gamma)}\right)+\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} e r r_{k} \tag{63}
\end{equation*}
$$

Now we can redefine the term $\operatorname{err}_{k}$ as $\operatorname{err}_{k}=\mathbb{E}_{s, a}\left(\left|A^{\pi_{\lambda_{k}}}-w^{k} \nabla \log \left(\pi^{\theta_{k}}(a \mid s)\right)\right|\right)$ and the above inequality will still hold as $\mathbb{E}(x) \leq \mathbb{E}(|x|)$ for any random variable $x$. Now as in Equation equation 22 we have

$$
\begin{align*}
\operatorname{err}_{k} & =\mathbb{E}_{s, a}\left(\left|A^{\pi_{\lambda_{k}}}-w^{k} \nabla \log \left(\pi^{\theta_{k}}(a \mid s)\right)\right|\right)  \tag{64}\\
& \leq \mathbb{E}_{s, a}\left(\left|A^{\pi_{\lambda_{k}}}-A_{k, J}\right|\right)+\mathbb{E}_{s, a}\left(\mid A_{k, J}-w^{k} \nabla \log \left(\pi^{\theta_{k}}(a \mid s) \mid\right)\right) \\
& \leq \underbrace{\mathbb{E}_{s, a}\left(\left|A^{\pi_{\lambda_{k}}}-A_{k, J}\right|\right)}_{I I}+\underbrace{\mathbb{E}_{s, a}\left(\left|A_{k, J}-w^{k} \nabla \log \left(\pi^{\theta_{k}}(a \mid s)\right)\right|\right)}_{I} \tag{65}
\end{align*}
$$

where $A_{k, j}$ is the estimate of $A^{\pi_{\lambda_{k}}}$ obtained at the $k^{t h}$ iteration of Algorithm 1 and $s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}$. We first derive bounds on $I$. From the definition of advantage function we have

$$
\begin{align*}
\mathbb{E}\left(\left|A^{\pi_{\lambda_{k}}}(s, a)-A_{k, J}(s, a)\right|\right)= & \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}} \mid Q^{\pi_{\lambda_{k}}}(s, a)-E_{a_{i+1} \sim \pi^{\lambda_{k}}} Q^{\pi_{\lambda_{k}}}\left(s, a_{i+1}\right) \\
& -Q_{k, J}(s, a)+E_{a_{i+1} \sim \pi^{\lambda_{k}}} Q_{k, J}\left(s, a_{i+1}\right) \mid  \tag{66}\\
= & \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}} \mid Q^{\pi_{\lambda_{k}}}(s, a)-E_{a_{i+1} \sim \pi^{\lambda_{k}}} Q^{\pi_{\lambda_{k}}}\left(s, a_{i+1}\right) \\
& -Q_{k, J}(s, a)+E_{a_{i+1} \sim \pi^{\lambda_{k}}} Q_{k, J}\left(s, a_{i+1}\right) \mid  \tag{67}\\
\leq & \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right| \\
+ & \left|\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a_{i+1} \sim \pi^{\lambda_{k}}}\right| Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a) \mid \tag{68}
\end{align*}
$$

We write the second term on the right hand side of Equation equation 68 as $\int\left(\mid\left(Q^{\pi_{\lambda_{k}}}(s, a)-\right.\right.$ $\left.Q_{k, J \mid}(s, a)\right) d\left(\mu_{k}\right)$ where $\mu_{k}$ is the measure associated with the state action distribution given by $s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{\lambda_{k}}$. Then we have

$$
\begin{equation*}
\int\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right| d\left(\mu_{k}\right) \leq\left\|\frac{d \mu_{k}}{d \mu^{*}}\right\|_{\infty} \int\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right| d\left(\mu^{*}\right) \tag{69}
\end{equation*}
$$

where $\mu^{*}$ is the measure associated with the state action distribution given by $s \sim d_{\nu}^{\pi^{*}}, a_{i+1} \sim \pi^{*}$. Now before we proceed further, we would like to introduce some notation for convenience, for two probability measures $\mu_{1}$ and $\mu_{2}$ on we define $\left\|\frac{d \mu_{1}}{d \mu_{2}}\right\|_{\infty}=\phi_{\mu_{1}, \mu_{2}}$

Thus Equation equation 69 becomes

$$
\begin{equation*}
\int\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right| d\left(\mu_{k}\right) \leq\left(\phi_{\mu_{k}, \mu^{*}}\right) \int\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right| d\left(\mu^{*}\right) \tag{70}
\end{equation*}
$$

Since $\int\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right| d\left(\mu^{*}\right)=\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a_{i+1} \sim \pi^{*}}\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right|$
Equation equation 68 now becomes.

$$
\begin{equation*}
\mathbb{E}\left|A^{\pi_{\lambda_{k}}}(s, a)-A_{k, J}(s, a)\right| \leq\left(1+\phi_{\mu_{k}, \mu^{*}}\right) \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right| \tag{71}
\end{equation*}
$$

Therefore minimizing $\left|A^{\pi_{\lambda_{k}}}(s, a)-A_{k, J}(s, a)\right|$ is equivalent to minimizing $\mid Q^{\pi_{\lambda_{k}}}(s, a)-$ $Q_{k, J}(s, a) \mid$.
In order to prove the bound on $\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left|Q^{\pi_{\lambda_{k}}}(s, a)-Q_{k, J}(s, a)\right|$ we first define some notation, let $Q_{1}, Q_{2}$ be two real valued functions on the state action space. The expression $Q_{1} \geq Q_{2}$ implies $Q_{1}(s, a) \geq Q_{2}(s, a) \forall(s, a) \in \mathcal{S} \times \mathcal{A}$.
Let $Q_{k, j}$ denotes our estimate of the action value function at iteration $k$ of Algorithm 1 and iteration $j$ of the first for loop of Algorithm 1. $Q^{\pi_{\lambda_{k}}}$ denotes the action value function induced by the policy $\pi_{\lambda_{k}}$.
Consider $\epsilon_{k, j+1}=T^{\pi_{\lambda_{k}}} Q_{k, j}-Q_{k, j+1}$.
Thus we get,

$$
\begin{align*}
Q^{\pi_{\lambda_{k}}}-Q_{k, j+1}= & T^{\pi_{\lambda_{k}}} Q^{\pi_{\lambda_{k}}}-Q_{k, j+1}  \tag{72}\\
= & T^{\pi_{\lambda_{k}}} Q^{\pi_{\lambda_{k}}}-T^{\pi_{\lambda_{k}}} Q_{k, j}+T^{\pi_{\lambda_{k}}} Q_{k, j} \\
& -Q_{k, j+1}  \tag{73}\\
= & \gamma\left(P^{\pi_{\lambda_{k}}}\left(Q^{\pi_{\lambda_{k}}}-Q_{k, j}\right)\right)+\epsilon_{k, j+1}  \tag{74}\\
\left|Q^{\pi_{\lambda_{k}}}-Q_{k, j+1}\right| \leq & \gamma\left(P^{\pi_{\lambda_{k}}}\left(\left|Q^{\pi_{\lambda_{k}}}-Q_{k, j}\right|\right)\right)+\left|\epsilon_{k, j+1}\right| \tag{75}
\end{align*}
$$

Right hand side of Equation equation 72 is obtained by writing $Q^{\pi_{\lambda_{k}}}=T^{\pi_{\lambda_{k}}} Q^{\pi_{\lambda_{k}}}$. This is because the function $Q^{\pi_{\lambda_{k}}}$ is a stationary point with respect to the operator $T^{\pi_{\lambda_{k}}}$. Equation equation 73 is obtained from equation 72 by adding and subtracting $T^{\pi^{\lambda_{k}}}$. We get equation 75 from equation 74 by taking the absolute value on both sides and applying the triangle inequality on the right hand side.
By recursion on $k$, we get,

$$
\begin{equation*}
\left|Q^{\pi_{\lambda_{k}}}-Q_{k, J}\right| \leq \sum_{j=0}^{J-1} \gamma^{J-j-1}\left(P^{\pi_{\lambda_{k}}}\right)^{J-j-1}\left|\epsilon_{k, j+1}\right|+\gamma^{J}\left(P^{\pi_{\lambda_{k}}}\right)^{J}\left(\left|Q^{\pi_{\lambda_{k}}}-Q_{0}\right|\right) \tag{76}
\end{equation*}
$$

From this we obtain

$$
\begin{align*}
\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left|Q^{\pi_{\lambda_{k}}}-Q_{k, J}\right| \leq & \sum_{k=0}^{J-1} \gamma^{J-j-1} \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left(\left(P^{\pi_{\lambda_{k}}}\right)^{K-J-1}\left|\epsilon_{k, j+1}\right|\right) \\
& +\gamma^{J} \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left(P^{\pi_{\lambda_{k}}}\right)^{J}\left(\left|Q^{\pi_{\lambda_{k}}}-Q_{0}\right|\right) \tag{77}
\end{align*}
$$

For a fixed $j$ consider the term $\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left(\left(P^{\pi_{\lambda_{k}}}\right)^{J-j-1}\left|\epsilon_{k, j+1}\right|\right)$. We then write

$$
\begin{align*}
& \mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi_{j}}\left(\left(P^{\pi_{\lambda_{k}}}\right)^{J-j-1}\left|\epsilon_{k, j+1}\right|\right) \leq\left\|\frac{d\left(P^{\pi_{\lambda_{k}}}\right)^{J-j-1} \mu_{j}}{d \mu_{k}^{\prime}}\right\|_{\infty} \int\left|\epsilon_{k, j+1}\right| d \mu_{k}^{\prime}  \tag{78}\\
& \leq\left(\phi_{\mu_{k}^{\prime}, \mu_{j}}\right) \mathbb{E}_{(s, a) \sim \zeta_{\nu}^{\pi}}{ }_{\infty}{ }^{\lambda_{k}}(s, a)  \tag{79}\\
&\left(\left|\epsilon_{k, j+1}\right|\right) \mid
\end{align*}
$$

Here $\mu_{j}$ is the measure associated with the state action distribution given by sampling from $s \sim$ $d_{\nu}^{\pi^{*}}, a_{i+1} \sim \pi^{*}$ and then applying the operator $P^{\pi_{\lambda_{k}}}, J-j-1$ times. $\mu_{j}$ is the measure associated with the steady state action distribution given by $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)$. Thus Equation equation 77 becomes

$$
\begin{equation*}
\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left|Q^{\pi_{\lambda_{k}}}-Q_{k, J}\right| \leq \sum_{k=0}^{J-1} \gamma^{J-j-1}\left(\phi_{\mu_{k}^{\prime}, \mu_{j}}\right) \mathbb{E}_{(s, a) \sim \zeta_{\nu}^{\pi \lambda_{k}}(s, a)}\left(\left|\epsilon_{k, j+1}\right|\right) \left\lvert\,+\gamma^{J}\left(\frac{R_{\max }}{1-\gamma}\right)\right. \tag{80}
\end{equation*}
$$

We get the second term on the right hand side by noting that $\left(Q^{\pi_{\lambda_{k}}}-Q_{0}\right) \leq \frac{R_{\max }}{1-\gamma}$. Now splitting $\epsilon_{k, j+1}$ as was done in Equation equation 25 we obtain

$$
\begin{align*}
\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left|Q^{\pi_{\lambda_{k}}}-Q_{k, J}\right| \leq & \sum_{j=0}^{J-1} \gamma^{J-j-1}\left(\left(\phi_{\mu_{k}^{\prime}, \mu_{j}}\right) \mathbb{E}\left|\epsilon_{k, j+1}^{1}\right|+\left(\phi_{\mu_{k}^{\prime}, \mu_{j}}\right) \mathbb{E}\left|\epsilon_{k, j+1}^{2}\right|\right. \\
& \left.+\left(\phi_{\mu_{k}^{\prime}, \mu_{j}}\right) \mathbb{E}\left|\epsilon_{k, j+1}^{3}\right|+\left(\phi_{\mu_{k}^{\prime}, \mu_{j}}\right) \mathbb{E}\left|\epsilon_{k, j+1}^{4}\right|\right)+\gamma^{J}\left(\frac{R_{\max }}{1-\gamma}\right) \tag{81}
\end{align*}
$$

Now using Lemmas $3,4,5,6$ we have for $m \geq \mathcal{O}\left(K . J . \delta^{-1}\right)$ with probability at-least $1-\delta$ we have

$$
\begin{align*}
\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left(Q^{\pi_{\lambda_{k}}}-Q_{k, j}\right) & \leq \sum_{j=0}^{J-1}\left(\mathcal{O}\left(\frac{1}{\sqrt{n}}\right)+\mathcal{O}\left(\sqrt{\epsilon_{\text {approx }}}\right)+\mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{n}\right. \\
& \left.+\mathcal{O}\left(\gamma^{J}\right)\right) \tag{82}
\end{align*}
$$

Note we had to increase the requirement of $m$ from $m \geq \mathcal{O}\left(\delta^{-1}\right)$ to $m \geq \mathcal{O}\left(K . J . \delta^{-1}\right)$ as we want the probability statement of lemma 6 to hold for all iterations of the outer for loop and the first inner for loop. Also note that for lemma 5 the number of iterations of gradient descent is $n$, so we replace $T_{k, j}$ to $n$.

From Equation equation 71 we get that for we have for $m \geq \mathcal{O}\left(K . J . \delta^{-1}\right)$ with probability at-least $1-\delta$ we have

$$
\begin{align*}
\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left(A^{\pi_{\lambda_{k}}}-A_{k, j}\right) & \leq \sum_{j=0}^{J-1}\left(\mathcal{O}\left(\frac{1}{\sqrt{n}}\right)+\mathcal{O}\left(\sqrt{\epsilon_{\text {approx }}}\right)+\mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{n}\right. \\
& \left.+\mathcal{O}\left(\gamma^{J}\right)\right) \tag{83}
\end{align*}
$$

We now derive bounds on $I I$. Note that $I I$ can be upper bounded as
$\mathbb{E}_{s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}}\left(\left|A_{k, J}-w^{k} \nabla \log \left(\pi^{\theta_{k}}(a \mid s)\right)\right| \leq\left(\phi_{\alpha_{k}, \mu^{*}}\right) \mathbb{E}_{s, a \sim \zeta_{\nu}^{\pi \lambda_{k}}}\left(\mid A_{k, J}-w^{k} \nabla \log \left(\pi^{\theta_{k}}(a \mid s) \mid\right)\right.\right.$
Here $\alpha_{k}$ is the measure corresponding to $s, a \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)$ and $\mu^{*}$ is as defined previously.
Now from lemma if we have $\beta_{i, k}=\frac{2}{\mu_{k}(i+1)}$, where $\mu_{k}$ is the strong convexity parameter for the loss function in $F_{k}(w)=\mathbb{E}_{s, a \sim \zeta_{\pi}^{\nu}}\left(\left|A_{k, J}-w^{k} \nabla \log \left(\pi^{\theta_{k}}(a \mid s)\right)\right|\right)^{2}$, we obtain from lemma 7 that

$$
\begin{equation*}
\left\|w_{k}-w_{k}^{*}\right\|_{2} \leq \mathcal{O}\left(\frac{\log (n)}{n}\right) \tag{85}
\end{equation*}
$$

Now define the function. From this definition and the fact that $w_{k}^{*}$ is also the minimizer of $I I$ we obtain

$$
\begin{equation*}
F_{k}\left(w_{k}\right)-F_{k}\left(w^{*}\right) \leq l_{F_{k}}\left\|w_{k}-w^{*}\right\|_{2} \leq \mathcal{O}\left(\frac{\log (n)}{n}\right) \tag{86}
\end{equation*}
$$

where $l_{F_{k}}$ is lipschitz constant of $F_{k}(w)$. Thus we obtain

$$
\begin{equation*}
F_{k}\left(w_{k}\right)-F_{k}\left(w^{*}\right) \leq \mathcal{O}\left(\frac{\log (n)}{n}\right) \tag{87}
\end{equation*}
$$

which gives us

$$
\begin{equation*}
F_{k}\left(w_{k}\right) \leq \mathcal{O}\left(\frac{\log (n)}{n}\right)+\epsilon_{b i a s} \tag{88}
\end{equation*}
$$

We get equation 88 from equation 87 by using assumption 4 .
Now from equation 84 we get

$$
\begin{equation*}
I I \leq \mathcal{O}\left(\frac{\log (n)}{n}\right)+\mathcal{O}\left(\epsilon_{b i a s}\right) \tag{89}
\end{equation*}
$$

Plugging Equations equation 88 and equation 83 in Equation equation 63 we get for $m \geq$ $\mathcal{O}\left(K . J . \delta^{-1}\right)$ with probability at-least $1-\delta$ we have

$$
\begin{align*}
\min _{k \leq K}\left(V^{*}(\nu)-V^{\pi_{\lambda_{K}}}(\nu)\right) & \leq \frac{1}{K} \sum_{k=1}^{K}\left(V^{*}(\nu)-V^{\pi_{\lambda_{K}}}(\nu)\right)  \tag{90}\\
\leq & \mathcal{O}\left(\frac{1}{\sqrt{K}(1-\gamma)}\right)+\frac{1}{K(1-\gamma)} \sum_{k=1}^{K}\left(\mathcal{O}\left(\frac{\log (J . n)}{J . n}\right)+\mathcal{O}\left(\gamma^{J}\right)\right) \\
& +\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} \sum_{j=0}^{J-1}\left(\mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{n}+\mathcal{O}\left(\frac{1}{\sqrt{n}}\right)\right) \\
& +\frac{1}{1-\gamma}\left(\mathcal{O}\left(\epsilon_{\text {bias }}\right)+\mathcal{O}\left(\sqrt{\epsilon_{\text {approx }}}\right)\right) \tag{91}
\end{align*}
$$

## D Proof of Supporting Lemmas

## D. 1 Proof Of Lemma 3

Proof. Using Assumption 3 and the definition of $Q_{k j_{1}}$ for some iteration $k$ of Algorithm 1 we have

$$
\begin{equation*}
\mathbb{E}_{s, a}\left(T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}\right)^{2} \leq \epsilon_{\text {approx }} \tag{92}
\end{equation*}
$$

where $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a)$.

Since $|a|^{2}=a^{2}$ we obtain

$$
\begin{equation*}
\mathbb{E}\left(\left|T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}\right|\right)^{2} \leq \epsilon_{\text {approx }} \tag{93}
\end{equation*}
$$

We have for a random variable $x, \operatorname{Var}(x)=\mathbb{E}\left(x^{2}\right)-(\mathbb{E}(x))^{2}$ hence $\mathbb{E}(x)=\sqrt{\mathbb{E}\left(x^{2}\right)-\operatorname{Var}(x)}$, Therefore replacing $x$ with $\left|T^{\pi_{\lambda_{k}}} Q^{\pi_{\lambda_{k}}}-Q_{k 1}\right|$ we get
using the definition of the variance of a random variable we get

$$
\begin{align*}
\mathbb{E}\left(\left|T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}\right|\right) & =\sqrt{\mathbb{E}\left(\left|T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}\right|\right)^{2}-\operatorname{Var}\left(\left|T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}\right|\right)}  \tag{94}\\
& \leq \sqrt{\mathbb{E}\left(\left|T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1}\right|\right)^{2}} \tag{95}
\end{align*}
$$

Therefore by definition of $Q_{k, j}^{1}$ and assumption 4 we have

$$
\begin{equation*}
\mathbb{E}\left(T^{\pi_{\lambda_{k}}} Q_{k, j-1}-Q_{k, j}^{1} \mid\right) \leq \sqrt{\epsilon_{\text {approx }}} \tag{96}
\end{equation*}
$$

Since $\epsilon_{k_{1}}=T^{\pi \lambda_{k}} Q^{\pi_{\lambda_{k}}}-Q_{k 1}$ we have

$$
\begin{equation*}
\mathbb{E}\left(\left|\epsilon_{k, j_{1}}\right|\right) \leq \sqrt{\epsilon_{\text {approx }}} \tag{97}
\end{equation*}
$$

## D. 2 Proof Of Lemma 4

Proof. From Lemma 2, we have

$$
\begin{equation*}
\underset{f_{\theta}}{\arg \min } \mathbb{E}_{x, y}\left(f_{\theta}(x)-g(x, y)\right)^{2}=\underset{f_{\theta}}{\arg \min }\left(\mathbb{E}_{x, y}\left(f_{\theta}(x)-\mathbb{E}\left(g\left(y^{\prime}, x\right) \mid x\right)\right)^{2}\right) \tag{98}
\end{equation*}
$$

We label $x$ to be the state action pair $(s, a), y$ is the next state action pair denoted by $\left(s^{\prime}, a_{i+1}\right)$. The function $f_{\theta}(x)$ to be $Q_{\theta}(s, a)$ and $g(x, y)$ to be the function $r^{\prime}(s, a)+\gamma Q_{k, j-1}\left(s^{\prime}, a_{i+1}\right)$
Then the loss function corresponding to Equation equation 36 becomes

$$
\begin{equation*}
\mathbb{E}\left(Q_{\theta}(s, a)-\left(r(s, a)+\gamma \mathbb{E} Q_{k, j-1}\left(s^{\prime}, a_{i+1}\right)\right)\right)^{2} \tag{99}
\end{equation*}
$$

where $(s, a) \sim \zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a), s^{\prime} \sim P(. \mid(s, a)), a_{i+1} \sim \pi^{\lambda_{k}}\left(. \mid s^{\prime}\right)$ and $r(s, a) \sim R(. \mid s, a)$.
Therefore by Lemma 2, we have that the function $Q_{\theta}(s, a)$ which minimizes Equation equation 99 it will be minimizing

$$
\begin{equation*}
\mathbb{E}_{s \sim d_{\nu}^{\pi_{\lambda}}, a \sim \pi_{\lambda_{k}}}\left(Q_{\theta}(s, a)-\mathbb{E}_{\left.s^{\prime} \sim P\left(s^{\prime} \mid s, a\right), r \sim \mathcal{R}(. \mid s, a)\right)}\left(r(s, a)+\gamma \mathbb{E} Q_{k, j-1}\left(s^{\prime}, a_{i+1}\right) \mid s, a\right)\right)^{2} \tag{100}
\end{equation*}
$$

But we have from Equation that

$$
\begin{equation*}
\mathbb{E}_{\left.s^{\prime} \sim P\left(s^{\prime} \mid s, a\right), r \sim R(\cdot \mid s, a)\right)}\left(r(s, a)+\gamma \mathbb{E} Q_{k, j-1}\left(s^{\prime}, a_{i+1}\right) \mid s, a\right)=T^{\pi_{\lambda_{k}}} Q_{k, j-1} \tag{101}
\end{equation*}
$$

Combining Equation equation 99 and equation 101 we get

$$
\begin{equation*}
\underset{Q_{\theta}}{\arg \min } \mathbb{E}\left(Q_{\theta}(s, a)-\left(r(s, a)+\gamma Q_{k, j-1}\left(s^{\prime}, a_{i+1}\right)\right)\right)^{2}=\underset{Q_{\theta}}{\arg \min } \mathbb{E}\left(Q_{\theta}(s, a)-T^{\pi_{\lambda_{k}}} Q_{k, j-1}\right)^{2} \tag{102}
\end{equation*}
$$

The left hand side of Equation equation 102 is $Q_{k, j}^{2}$ as defined in Definition 2 and the right hand side is $Q_{k, j}^{1}$ as defined in Definition 1, which gives us

$$
\begin{equation*}
Q_{k, j}^{2}=Q_{k, j}^{1} \tag{103}
\end{equation*}
$$

## D. 3 Proof Of Lemma 5

Proof. We define $R_{X_{k, j}, Q_{k, j-1}}(\theta)$ as

$$
R_{X_{k, j}, Q_{k, j-1}}(\theta)=\frac{1}{n} \sum_{\left(s_{i}, a_{i}\right) \in X_{k, j}}\left(Q_{\theta}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j-1}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2}
$$

Here, where $s, a$ are sampled from a Markov chain whose stationary distribution is, $(s, a) \sim$ $\zeta_{\nu}^{\pi_{\lambda_{k}}}(s, a) . Q_{\theta}$ is is the neural network corresponding to the parameter $\theta$ and $Q_{k, j-1}$ is the estimate of the $Q$ function obtained at iteration $k$ of the outer for loop and iteration $j-1$ of the first inner for loop of Algorithm 1.

We also define the term

$$
\begin{equation*}
L_{Q_{k, j-1}}\left(Q_{\theta}\right)=\mathbb{E}\left(Q_{\theta}(s, a)-\left(r^{\prime}(s, a)+\gamma Q_{k, j-1}\left(s^{\prime}, a^{\prime}\right)\right)^{2}\right. \tag{104}
\end{equation*}
$$

where $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}, r^{\prime}(\cdot \mid s, a) \sim R(\cdot \mid s, a), a_{i+1} \sim \pi_{\lambda_{k}}$

We denote by $\theta_{k, j}^{2}, \theta_{k, j}^{3}$ the parameters of the neural networks $Q_{k, j}^{2}, Q_{k, j}^{3}$ respectively. $Q_{k, j}^{2}, Q_{k, j}^{3}$ are defined in Definition 2 and 3 respectively.

We then obtain,

$$
\left.\left.\begin{array}{rl}
R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)-R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right) \leq & R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)-R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right) \\
& +L_{Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)-L_{Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)
\end{array}\right\} \begin{array}{rl}
R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)-L_{Q_{k, j-1}}\left(\theta_{k, j}^{2}\right) \\
& -R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)+L_{Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)
\end{array}\right\} \begin{array}{ll}
\mid R_{X_{k, j}, Q_{k, j-1}\left(\theta_{k, j}^{2}\right)-L_{Q_{k, j-1}}\left(\theta_{k, j}^{2}\right) \mid}^{I} \\
& +\underbrace{\left|R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)-L_{Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)\right|}_{I I}
\end{array}
$$

We get the inequality in Equation equation 105 because $L_{Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)-L_{Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)>0$ as $Q_{k, j}^{2}$ is the minimizer of the loss function $L_{Q_{k, j-1}}\left(Q_{\theta}\right)$.
Consider Lemma 1. The loss function $R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)$ can be written as the mean of loss functions of the form $l\left(a_{\theta}\left(s_{i}, a_{i}, s_{i+1}, a_{i+1}\right), y_{i}\right)$ where $l$ is the square function. $a_{\theta}\left(s_{i}, a_{i}, s_{i+1}, a_{i+1}\right)=$ $Q_{\theta}\left(s_{i}, a_{i}\right)$ and $y_{i}=\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j-1}\left(s_{i+1}, a_{i+1}\right)\right)$. Thus we have

$$
\begin{gather*}
\mathbb{E} \sup _{\theta \in \Theta^{\prime}}\left|R_{X_{k, j}, Q_{k, j-1}}(\theta)-L_{Q_{k, j-1}}(\theta)\right| \leq  \tag{108}\\
2 \eta^{\prime} \mathbb{E}\left(\operatorname{Rad}\left(\mathcal{A} \circ\left\{\left(s_{1}, a_{1}, s_{2}, a_{2}\right),\left(s_{2}, a_{2}, s_{3}, a_{3}\right), \cdots,\left(s_{n-1}, a_{n-1}, s_{n}, a_{n}\right)\right\}\right)\right)
\end{gather*}
$$

Note that the expectation is over all $\left(s_{i}, a_{i}\right)$ and that the set $\Theta^{\prime}=\left\{Q_{\theta_{k, j}^{2}}, Q_{\theta_{k, j}^{3}}\right\}$. We use this set because we only need this inequality to hold for $Q_{\theta_{k, j}^{2}}$ and $Q_{\theta_{k, j}^{3}}$. Where $n=\left|X_{k, j}\right|,(\mathcal{A} \circ$ $\left\{\left(s_{1}, a_{1}\right),\left(s_{2}, a_{2}\right),\left(s_{3}, a_{3}\right), \cdots,\left(s_{n}, a_{n}\right)\right\}=\left\{Q_{\theta}\left(s_{1}, a_{1}\right), Q_{\theta}\left(s_{2}, a_{2}\right), \cdots, Q_{\theta}\left(s_{n}, a_{n}\right)\right\}$ and $\eta_{i+1}$ is the Lipschitz constant for the square function over the state action space $[0,1]^{d}$. The expectation is with respect to $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}, s_{i}^{\prime} \sim P\left(s^{\prime} \mid s, a\right) r_{i} \sim R\left(. \mid s_{i}, a_{i}\right)_{i \in(1, \cdots, n)}$,
Now from theorem 5 and theorem 1 of Bertail \& Portier (2019) we have that

$$
\begin{equation*}
\left(\operatorname{Rad}\left(\mathcal{A} \circ\left\{\left(s_{1}, a_{1}, s_{2}, a_{2}\right),\left(s_{2}, a_{2}, s_{3}, a_{3}\right), \cdots,\left(s_{n-1}, a_{n-1}, s_{n}, a_{n}\right)\right\}\right)\right) \leq C_{k} \frac{1}{\sqrt{n}} \tag{109}
\end{equation*}
$$

Note that in Bertail \& Portier (2019) a factor of $\log \log (n)$ in the numerator is introduced in later theorems, we ignore that factor due to the fact that it is practically constant and the we will have a factor of $\log (n)$ from the error incurred in the actor step.
We use this result as the state action pairs are drawn not from the stationary state of the policy $\pi_{\lambda_{k}}$ but from a Markov chain with the same steady state distribution. Thus we have

$$
\begin{equation*}
\mathbb{E}\left|\left(R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)\right)-L_{Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)\right| \leq C_{k} \frac{1}{\sqrt{n}} \tag{110}
\end{equation*}
$$

The same argument can be applied for $Q_{k, j}^{3}$ to get

$$
\begin{equation*}
\mathbb{E}\left|\left(R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)\right)-L_{Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)\right| \leq C_{k} \frac{1}{\sqrt{n}} \tag{111}
\end{equation*}
$$

Then we have

$$
\begin{equation*}
\mathbb{E}\left(R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{2}\right)-R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)\right) \leq C_{k} \frac{1}{\sqrt{n}} \tag{112}
\end{equation*}
$$

Plugging in the definition of $R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{2}\right), R_{X_{k, j}, Q_{k, j-1}}\left(\theta_{k, j}^{3}\right)$ in equation equation 112 and denoting $C_{k} \frac{1}{\sqrt{n}}$ as $\epsilon$ we get

$$
\begin{align*}
& \frac{1}{n} \sum_{i=1}^{n}\left(\mathbb{E}\left(Q_{k, j}^{2}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{2}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2}\right. \\
& \left.-\mathbb{E}\left(Q_{k, j}^{3}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{3}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2}\right) \leq \epsilon \tag{113}
\end{align*}
$$

Now for a fixed $i$ consider the term $\alpha_{i}$ defined as.

$$
\begin{array}{r}
\mathbb{E}_{s_{i+1} \sim P\left(. \mid s_{i}, a_{i}\right)}\left(Q_{k, j}^{2}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{2}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2} \\
-\mathbb{E}_{s_{i+1} \sim P\left(. \mid s_{i}, a_{i}\right)}\left(Q_{k, j}^{3}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{3}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2} \tag{114}
\end{array}
$$

where $s_{i}, a_{i}, s_{i+1}, a_{i+1}$ are drawn from the state action distribution at the $i^{t h}$ step of the Markov chain induced by following the policy $\pi_{\lambda_{k}}$.
Now for a fixed $i$ consider the term $\beta_{i}$ defined as.

$$
\begin{array}{r}
\mathbb{E}_{s_{i+1} \sim P\left(. \mid s_{i}, a_{i}\right)}\left(Q_{k, j}^{2}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{2}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2} \\
-\mathbb{E}_{s_{i+1} \sim P\left(. \mid s_{i}, a_{i}\right)}\left(Q_{k, j}^{3}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{3}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2} \tag{115}
\end{array}
$$

where $s_{i}, a_{i}, s_{i+1}, a_{i+1}$ are drawn from the steady state action distribution with $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}$. Note here that $\alpha_{i}$ and $\beta_{i}$ are the same function with only the state action pairs being drawn from different distributions.

Using these definitions we obtain

$$
\begin{align*}
\left|\mathbb{E}\left(\alpha_{i}\right)-\mathbb{E}\left(\beta_{i}\right)\right| & \leq \sup _{\left(s_{i}, a_{i}\right)}\left|2 \cdot \max \left(\alpha_{i}, \beta_{i}\right)\right|\left(\kappa_{i}\right)  \tag{116}\\
& \leq\left(4 \frac{R}{1-\gamma}\right)^{2} p \rho^{i} \tag{117}
\end{align*}
$$

We obtain Equation equation 116 by using the identity $\left|\int f d \mu-\int f d \nu\right| \leq$ $\left|\max _{\mathcal{S} \times \mathcal{A}}(f)\right| \sup _{\mathcal{S} \times \mathcal{A}} \int(d \mu-d \nu)\left|\leq\left|\max _{\mathcal{S} \times \mathcal{A}}(f)\right| d_{T V}(\mu, \nu)\right|$, where $\mu$ and $\nu$ are two $\sigma$ finite state action probability measures and $f$ is a bounded measurable function. We have used $\kappa_{i}$ to represent the total variation distance between the state action measures of the steady state action distribution denoted by $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}$ and the state action distribution at the $i^{t h}$ step of the Markov chain induced by following the policy $\pi^{\lambda_{k}}$. The expectation is with respect to $\left(s_{i}, a_{i}\right)$. We obtain Equation equation 117 from Equation equation 116 by using Assumption 2 and the fact that $\alpha_{i}$ and $\beta_{i}$ are upper bounded by $\left(4 \frac{R}{1-\gamma}\right)^{2}$
From equation equation 117 we get

$$
\begin{equation*}
\mathbb{E}\left(\alpha_{i}\right) \geq \mathbb{E}\left(\beta_{i}\right)-4\left(\frac{R}{1-\gamma}\right)^{2} p \rho^{i} \tag{118}
\end{equation*}
$$

We get Equation equation 118 from Equation equation 117 using the fact that $|a-b| \leq c$ implies that $(-c \geq(a-b) \leq c)$ which in turn implies $a \geq b-c$.

Using Equation equation 118 in equation equation 115 we get

$$
\begin{align*}
& \frac{1}{n} \sum_{i=1}^{n}\left(\mathbb{E}\left(Q_{k, j}^{2}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{2}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2}\right. \\
& \left.-\mathbb{E}\left(Q_{k, j}^{3}\left(s_{i}, a_{i}\right)-\left(r^{\prime}\left(s_{i}, a_{i}\right)+\gamma Q_{k, j}^{3}\left(s_{i+1}, a_{i+1}\right)\right)\right)^{2}\right) \\
\leq & \epsilon+\frac{1}{n} \sum_{i=1}^{n} 4\left(\frac{R}{1-\gamma}\right)^{2} p \rho^{i} \\
\leq & \epsilon+\frac{1}{n} 4\left(\frac{R}{1-\gamma}\right)^{2} p \frac{1}{1-\rho} \tag{119}
\end{align*}
$$

In Equation equation $119\left(s_{i}, a_{i}\right)$ are now drawn from $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}$ for all $i$.
We ignore the second term on the right hand side as it is $\tilde{\mathcal{O}}\left(\frac{1}{n}\right)$ as compared to the first term which is $\tilde{\mathcal{O}}\left(\frac{1}{\sqrt{n}}\right)$. Additionally the expectation in Equation equation 119 is with respect to $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}, r^{\prime}(\cdot \mid s, a) \sim R(\cdot \mid s, a), s_{i+1} \sim P\left(. \mid s_{i}, a_{i}\right), a_{i+1} \sim \pi^{\lambda_{k}}$

Since now we have $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}$ for all $i$, Equation equation 119 is equivalent to,

$$
\begin{equation*}
\mathbb{E} \underbrace{\left(Q_{k, j}^{2}(s, a)-Q_{k, j}^{3}(s, a)\right)}_{A 1} \underbrace{\left(Q_{k, j}^{2}(s, a)+Q_{k, j}^{3}(s, a)-2\left(r^{\prime}(s, a)\right)+\gamma Q_{k, j-1}\left(s^{\prime}, a_{i+1}\right)\right)}_{A 2} \leq \epsilon \tag{120}
\end{equation*}
$$

Where the expectation is now over $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}, r^{\prime}(s, a) \sim R(. \mid s, a)$ and $s^{\prime} \sim P(. \mid s, a), a_{i+1} \sim$ $\pi_{\lambda_{k}}$. We re-write Equation equation 120 as

$$
\begin{align*}
\int \underbrace{\left(Q_{k, j}^{2}(s, a)-Q_{k, j}^{3}(s, a)\right)}_{A 1} & \times \\
& \times \underbrace{\left(Q_{k, j}^{2}(s, a)+Q_{k, j}^{3}(s, a)-2\left(r^{\prime}(s, a)\right)+\gamma \max _{a \in \mathcal{A}} Q_{k, j-1}\left(s^{\prime}, a\right)\right)}_{A 2} \times \\
& \times d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right) \leq \epsilon \tag{121}
\end{align*}
$$

Where $\mu_{1}$ is the state action distribution $(s, a) \sim \zeta_{\pi_{\lambda_{k}}}^{\nu}, \mu_{2}, \mu_{3}, \mu_{4}$ are the measures with respect to $(s, a), r^{\prime}, s^{\prime}$ and $a_{i+1}$ respectively.
Now for the integral in Equation equation 121 we split the integral into four different integrals. Each integral is over the set of $(s, a), r^{\prime}, s^{\prime}, a_{i+1}$ corresponding to the 4 different combinations of signs of $A 1, A 2$.

$$
\begin{align*}
& \int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1 \geq 0, A 2 \geq 0}(A 1)(A 2) d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right) \\
& +\int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1<0, A 2<0}(A 1)(A 2) d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right) \\
& +\int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1 \geq 0, A 2<0}(A 1)(A 2) d \mu_{1}(s, a) d \mu_{2}(r) d \nu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right) \\
& +\int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1<0, A 2 \geq 0}(A 1)(A 2) d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right) \leq \epsilon \tag{122}
\end{align*}
$$

Now note that the first 2 terms are non-negative and the last two terms are non-positive. We then write the first two terms as

$$
\begin{align*}
& \int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1 \geq 0, A 2 \geq 0}(A 1)(A 2) d(s, a) d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right) \\
= & C_{k, j_{1}} \int\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right| d \mu_{1} \\
= & C_{k, j_{1}} \mathbb{E}\left(\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right|\right)_{\mu_{1}}  \tag{123}\\
& \int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1<0, A 2<0}(A 1)(A 2) d(s, a) d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right) \\
= & C_{k, j_{2}} \int\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right| d \nu \\
= & C_{k, j_{2}} \mathbb{E}\left(\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right|\right)_{\mu_{1}} \tag{124}
\end{align*}
$$

We write the last two terms as

$$
\begin{align*}
& \int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1 \geq 0, A 2<0}(A 1)(A 2) d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right)=C_{k, j_{3}} \epsilon  \tag{125}\\
& \int_{\left\{(s, a), r^{\prime}, s^{\prime}\right\}: A 1<0, A 2 \geq 0}(A 1)(A 2) d \mu_{1}(s, a) d \mu_{2}(r) d \mu_{3}\left(s^{\prime}\right) d \mu_{4}\left(a_{i+1}\right)=C_{k, j_{4}} \epsilon \tag{126}
\end{align*}
$$

Here $C_{k, j_{1}}, C_{k, j_{2}}, C_{k, j_{4}}$ and $C_{k, j_{4}}$ are positive constants. Plugging Equations equation 123, equation 124, equation 125, equation 126 into Equation equation 121.

$$
\begin{equation*}
\left(C_{k, j_{1}}+C_{k, j_{2}}\right) \mathbb{E}\left(\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right|\right)_{\mu_{1}}-\left(C_{k, j_{3}}+C_{k, j_{4}}\right) \epsilon \leq \epsilon \tag{127}
\end{equation*}
$$

which implies

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right|\right)_{\mu_{1}} \leq\left(\frac{1+C_{k, j_{3}}+C_{k, j_{4}}}{C_{k, j_{1}}+C_{k, j_{2}}}\right) \epsilon \tag{129}
\end{equation*}
$$

Thus we have

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right|\right)_{\mu_{1}} \leq\left(\frac{1+C_{k, j_{3}}+C_{k, j_{4}}}{C_{k, j_{1}}+C_{k, j_{2}}}\right) C_{k} \frac{1}{\sqrt{n}} \tag{131}
\end{equation*}
$$

which implies

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{k, j}^{2}-Q_{k, j}^{3}\right|\right)_{\mu_{1}} \leq \tilde{\mathcal{O}}\left(\frac{1}{\sqrt{n}}\right) \tag{133}
\end{equation*}
$$

## D. 4 Proof Of Lemma 6

Proof. Note that from theorem 5 of Allen-Zhu et al. (2019) we have that for a neural network with at least $m$ neurons in every layer, sample size of $n$, we have after $T_{k, j}$ steps of the gradient descent algorithm with step size $\alpha$ we have with probability at least $1-\exp (-\Omega(\log (m)))$

$$
\begin{equation*}
F\left(\theta^{T}\right) \leq\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{T_{k, j}} F\left(\theta^{0}\right) \tag{135}
\end{equation*}
$$

Where $F$ is the loss function to be minimized. The specific form of $F\left(\theta^{T}\right)$ is obtained in the proof of theorem 1 in the appendix. Note that we have ignored the minimum euclidean distance between training samples denoted as $\delta$ in Allen-Zhu et al. (2019). It is to be noted that this is non-zero for our case as for any two separate $s, a$ pairs the target function will be different for a non-zero neural network.

Now plug gin in $m \geq \mathcal{O}\left(\delta^{-1}\right)$ we obtain with probability at least $1-\delta$

$$
\begin{equation*}
F\left(\theta^{T_{k, j}}\right) \leq \mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{T_{k, j}} \tag{136}
\end{equation*}
$$

Now choose a constant $C_{k, j}^{\prime}$ such that

$$
\begin{equation*}
C_{k, j}^{\prime}\left|\theta^{T_{k, j}}-\theta^{*}\right| \leq F\left(\theta^{T}\right) \leq \mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{T_{k, j}} \tag{137}
\end{equation*}
$$

where $\theta^{*}$ is the global minima for the loss function $F$. Now using the Lipschitz property of neural networks with respect to the parameters as discussed in Reddi et al. (2019) we have
$\left|Q_{\theta^{T}}(s, a)-Q_{\theta^{*}}(s, a)\right| \leq\left(L_{k, j}\right)\left(C_{k, j}^{\prime}\right)\left(\theta^{T}-\theta^{*}\right) \leq\left(L_{k, j}\right) F\left(\theta^{T}\right) \leq \mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{T_{k, j}}$
for all $(s, a) \in \mathcal{S} \times \mathcal{A}$. Note that by definition $Q_{\theta^{T}}(s, a)$ is $Q_{k, j}$ and $Q_{\theta^{*}}(s, a)$ is $Q_{k, j}^{3}$, thus we have with probability at least $1-\delta$

$$
\begin{equation*}
\left|Q_{\theta^{T}}(s, a)-Q_{\theta^{*}}(s, a)\right| \leq \mathcal{O}\left(1-\Omega\left(\frac{\alpha m}{n^{2}}\right)\right)^{T_{k, j}} \tag{139}
\end{equation*}
$$

Now taking expectation on both side with respect to $(s, a) \sim \zeta_{\pi^{\lambda} k}^{\nu}$ gives us the required result.

## D. 5 Proof Of LEmma 7

Proof. Note that this lemma is a direct application of theorem 1 from Doan (2022). Note that in the second inner for loop of Algorithm 1 we perform an iterations for every state action pair sampled in the first inner for loop for a total of $J \cdot n$ samples and noting that we are taking a steps of gradient descent for each sample gives us the required result.

## E Comparison of Sample Complexity Analysis with Natural Policy Gradient

For natural policy gradient (NPG) (Agarwal et al., 2021), to derive the sample complexity result, the average error in estimation till iteration $K$ is given by

$$
\begin{equation*}
\underset{k \in\{1, \cdots, K\}}{\arg \min } V^{*}(\nu)-V^{\pi_{K}}(\nu) \leq\left(\frac{\log (|\mathcal{A}|)}{K \eta(1-\gamma)}+\frac{\eta \beta_{k} W^{2}}{2(1-\gamma)}+\frac{1}{K} \sum_{i=k}^{K} \frac{e r r_{k}}{1-\gamma}\right) \tag{140}
\end{equation*}
$$

where $e r r_{k}$ in the last term on the right-hand side of equation 140 is

$$
\begin{equation*}
\operatorname{err}_{k}=\mathbb{E}_{s, a}\left(A^{\pi_{\lambda_{K}}}-w^{k} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right) \tag{141}
\end{equation*}
$$

where $s \sim d_{\nu}^{\pi^{*}}, a \sim \pi^{*}(. \mid s), w_{k}$ is our estimate of the NPG gradient update term and $\lambda_{k}$ is the policy parameter.

The term $\operatorname{err}_{k}$ is then decomposed in the following manner

$$
\begin{align*}
\mathbb{E}_{s, a}\left(A^{\pi_{\lambda_{K}}}-w^{k} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right)=\mathbb{E}_{s, a} & \left(A^{\pi_{\lambda_{K}}}-w^{*} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right) \\
+ & \left(w^{*}-w^{k)} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right) \tag{142}
\end{align*}
$$

where $w^{*}=\arg \min _{w} \mathbb{E}_{s, a}\left(Q^{\pi_{\lambda_{K}}}-w \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right)$ where $s \sim d_{\nu}^{\pi^{\lambda_{k}}}, a \sim \pi^{\lambda_{k}}(. \mid s)$.
For ease of notation we define

$$
\begin{equation*}
\mathbb{E}_{s \sim d_{\nu}^{\pi}, a \sim \pi(. \mid s)}\left(Q^{\pi_{\lambda}}-w \nabla \log \left(\pi_{\lambda}(a \mid s)\right)\right)^{2}=L\left(w, \lambda, d_{\nu}^{\pi}\right) \tag{143}
\end{equation*}
$$

Equation equation 142 is then be upper bounded as

$$
\begin{align*}
\mathbb{E}_{s, a}\left(Q^{\pi_{\lambda_{K}}}-w^{k} \nabla \log \left(\pi_{\lambda_{k}}(a \mid s)\right)\right) \leq & \sqrt{L\left(w^{*}, \lambda_{k}, d_{\nu}^{\pi^{*}}\right)} \\
& +\phi_{k} \sqrt{L\left(w^{k}, \lambda_{k}, d_{\nu}^{\pi_{\lambda_{k}}}\right)-L\left(w^{*}, \lambda_{k}, d_{\nu}^{\pi_{\lambda_{k}}}\right)} \tag{144}
\end{align*}
$$

where $\phi_{k}$ is a constant which represents the change in expectation from $d_{\nu}^{\pi^{*}}$ to $d_{\nu}^{\pi_{\lambda_{k}}}$.
Assumption 6.1 and 6.2 in (Agarwal et al., 2021) are as follows

$$
\begin{align*}
L\left(w^{k}, \lambda_{k}, d_{\nu}^{\pi_{\lambda}}\right)-L\left(w^{*}, \lambda_{k}, d_{\nu}^{\pi_{\lambda_{k}}}\right) & \leq \epsilon_{\text {stat }}  \tag{145}\\
L\left(w^{*} ; \lambda_{k}, d^{\pi_{\nu}^{*}}\right) & \leq \epsilon_{\text {bias }} \tag{146}
\end{align*}
$$

$\forall k \in\{1, \cdots, K\}$, where $K$ is the total number of iterations of the NPG algorithm.
The assumption in Equation equation 145 is known as the excess risk assumption and places an upper bound on the error incurred due to the difference between the obtained estimate $w^{k}$ and the optimal solution $w^{*}$ which minimizes $L\left(w, \lambda_{k}, d_{\nu}^{\pi \lambda_{k}}\right)$. It is a measure of uncertainty in estimating the natural gradient update.
The assumption in Equation equation 146 is known as the transfer error assumption and places an upper bound on the loss function $L\left(w, \lambda_{k}, d_{\nu}^{\pi^{*}}\right)$ evaluated at the minima of the loss function $L\left(w ; \lambda^{k}, d_{\nu}^{\pi_{\lambda_{k}}}\right)$. This is a measure of how similar the policy $\pi^{\lambda_{k}}$ is to the optimal policy $\pi^{*}$.

In the analysis of Agarwal et al. (2021), using results of stochastic gradient descent on a convex loss function, $\epsilon_{\text {stat }}$ is assumed to be upper bounded as $\tilde{\mathcal{O}}\left(\frac{1}{\sqrt{n_{k}}}\right)$ where $n_{k}$ is the number of state action samples at iteration $k$. Further, $\epsilon_{\text {bias }}$ is directly assumed as a constant while it depends on the accurate estimation of $A^{\pi_{\lambda_{k}}}$.
Comparison. We note that the analysis in Equation equation 145-equation 146 does not consider (i) the extra $\left(\frac{1}{1-\gamma}\right)$ state action samples required to obtain Monte Carlo estimate $A^{\pi_{\lambda_{k}}}$. This is because each Monte Carlo estimate of $A^{\pi_{\lambda_{k}}}$ requires on average $\left(\frac{1}{1-\gamma}\right)$ state action samples; (ii) the error incurred due to gap between the Monte Carlo estimate $A^{\pi_{\lambda_{k}}}$ and the actual Q-function. In Agarwal et al. (2021), Monte Carlo estimate is only shown to be an unbiased estimate of $A^{\pi_{\lambda_{k}}}$ and no error bound for the estimate is given. This error bound will require additional samples to be very close such that the obtained value function for the policy is $\epsilon$-close. This is the key gap due to which our algorithm gets additional $1 / \epsilon$ in the sample complexity.

Our analysis considers the number of samples required to estimate $A^{\pi_{\lambda_{k}}}$ to a given accuracy in our sample complexity analysis. In order to account for the difference between the optimal policy $\pi^{*}$ and the policy estimate $\pi_{\lambda_{k}}$, which has been used and verified in prior works such as Farahmand et al. (2010).

The authors in (Liu et al., 2020b) also perform a similar analysis but only has an Assumption similar to Equation equation 146. This assumption also suffers from the same drawback described above.

## F Sample Complexity using Convex Reformulation with Two-Layer Neural Networks

We now demonstrate how the error bound in Theorem 1 can be made deterministic if we use a 2 layer neural network with ReLU activation functions.

This part of the Appendix will first go into detail as to how a 2 layer neural network can be reformulated as a convex problem. We then prove the supplementary lemmas and additional assumptions required to prove a deterministic version of Theorem 1.
A 2-layer ReLU Neural Network with input $x \in \mathbb{R}^{d}$ is defined as $f(x)=\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i}$, where $m \geq 1$ is the number of neurons in the neural network, the parameter space is $\Theta_{m}=\mathbb{R}^{d \times m} \times \mathbb{R}^{m}$ and $\theta=(U, \alpha)$ is an element of the parameter space, where $u_{i}$ is the $i^{t h}$ column of $U$, and $\alpha_{i}$ is the $i^{t h}$ coefficient of $\alpha$. The function $\sigma^{\prime}: \mathbb{R} \rightarrow \mathbb{R}_{\geq 0}$ is the ReLU or restricted linear function defined as $\sigma^{\prime}(x) \triangleq \max (x, 0)$. In order to obtain parameter $\theta$ for a given set of data $X \in \mathbb{R}^{n \times d}$ and the corresponding response values $y \in \mathbb{R}^{n \times 1}$, we desire the parameter that minimizes the squared loss, given by

$$
\begin{equation*}
\mathcal{L}(\theta)=\left\|\sum_{i=1}^{m} \sigma\left(X u_{i}\right) \alpha_{i}-y\right\|_{2}^{2} . \tag{147}
\end{equation*}
$$

In equation 147, we have the term $\sigma\left(X u_{i}\right)$ which is a vector $\left\{\sigma^{\prime}\left(\left(x_{j}\right)^{T} u_{i}\right)\right\}_{j \in\{1, \cdots, n\}}$, where $x_{j}$ is the $j^{\text {th }}$ row of $X$. It is the ReLU function applied to each element of the vector $X u_{i}$. We note that the optimization in Equation equation 147 is non-convex in $\theta$ due to the presence of the ReLU activation function. In Wang et al. (2021b), it is shown that this optimization problem has an equivalent convex form, provided that the number of neurons $m$ goes above a certain threshold value. This convex problem is obtained by replacing the ReLU functions in the optimization problem with equivalent diagonal operators. The convex problem is given as

$$
\begin{equation*}
\mathcal{L}_{\beta}^{\prime}(p):=\left\|\sum_{D_{i} \in D_{X}} D_{i}\left(X p_{i}\right)-y\right\|_{2}^{2}, \tag{148}
\end{equation*}
$$

where $p \in \mathbb{R}^{d \times\left|D_{X}\right|} . D_{X}$ is the set of diagonal matrices $D_{i}$ which depend on the data-set $X$. Except for cases of $X$ being low rank, it is not computationally feasible to obtain the set $D_{X}$. We instead use $\tilde{D} \in D_{X}$ to solve the convex problem in equation 148 where $p$ now would lie in $p \in \mathbb{R}^{d \times|\tilde{D}|}$.

For a set of parameters $\theta=(u, \alpha) \in \Theta$, we denote neural network represented by these parameters as

$$
\begin{equation*}
Q_{\theta}(s, a)=\sum_{i=1}^{m} \sigma^{\prime}\left((s, a)^{T} u_{i}\right) \alpha_{i} \tag{149}
\end{equation*}
$$

For representing the action value function, we will use a 2 layer ReLU neural network. In this section, we first lay out the theory behind the convex formulation of the 2 layer ReLU neural network. In the next section it will shown how it is utilised for the FQI algorithm.
In order to obtain parameter $\theta$ for a given set of data $X \in \mathbb{R}^{n \times d}$ and the corresponding response values $y \in \mathbb{R}^{n \times 1}$, we desire the parameter that minimizes the squared loss (with a regularization parameter $\beta \in[0,1]$ ), given by

$$
\begin{equation*}
\mathcal{L}(\theta)=\left\|\sum_{i=1}^{m} \sigma\left(X u_{i}\right) \alpha_{i}-y\right\|_{2}^{2} \tag{150}
\end{equation*}
$$

Here, we have the term $\sigma\left(X u_{i}\right)$ which is a vector $\left\{\sigma^{\prime}\left(\left(x_{j}\right)^{T} u_{i}\right)\right\}_{j \in\{1, \cdots, n\}}$ where $x_{j}$ is the $j^{\text {th }}$ row of $X$. It is the ReLU function applied to each element of the vector $X u_{i}$. We note that the optimization in Equation equation 147 is non-convex in $\theta$ due to the presence of the ReLU activation function. In Wang et al. (2021b), it is shown that this optimization problem has an equivalent convex form, provided that the number of neurons $m$ goes above a certain threshold value. This convex problem is obtained by replacing the ReLU functions in the optimization problem with equivalent diagonal operators. The convex problem is given as

$$
\begin{equation*}
\mathcal{L}_{\beta}^{\prime}(p):=\left\|\sum_{D_{i} \in D_{X}} D_{i}\left(X p_{i}\right)-y\right\|_{2}^{2} \tag{151}
\end{equation*}
$$

where $p \in \mathbb{R}^{d \times\left|D_{X}\right|} . D_{X}$ is the set of diagonal matrices $D_{i}$ which depend on the data-set $X$. Except for cases of $X$ being low rank it is not computationally feasible to obtain the set $D_{X}$. We instead use $\tilde{D} \in D_{X}$ to solve the convex problem

$$
\begin{equation*}
\mathcal{L}_{\beta}^{\prime}(p):=\left\|\sum_{D_{i} \in \tilde{D}} D_{i}\left(X p_{i}\right)-y\right\|_{2}^{2}, \tag{152}
\end{equation*}
$$

where $p \in \mathbb{R}^{d \times|\tilde{D}|}$. In order to understand the convex reformulation of the squared loss optimization problem, consider the vector $\sigma\left(X u_{i}\right)$

$$
\sigma\left(X u_{i}\right)=\left[\begin{array}{c}
\left\{\sigma^{\prime}\left(\left(x_{1}\right)^{T} u_{i}\right)\right\}  \tag{153}\\
\left\{\sigma^{\prime}\left(\left(x_{2}\right)^{T} u_{i}\right)\right\} \\
\vdots \\
\left\{\sigma^{\prime}\left(\left(x_{n}\right)^{T} u_{i}\right)\right\} .
\end{array}\right]
$$

Now for a fixed $X \in \mathbb{R}^{n \times d}$, different $u_{i} \in \mathbb{R}^{d \times 1}$ will have different components of $\sigma\left(X u_{i}\right)$ that are non zero. For example, if we take the set of all $u_{i}$ such that only the first element of $\sigma\left(X u_{i}\right)$ are non zero (i.e, only $\left(x_{1}\right)^{T} u_{i} \geq 0$ and $\left.\left(x_{j}\right)^{T} u_{i}<0 \forall j \in[2, \cdots, n]\right)$ and denote it by the set $\mathcal{K}_{1}$, then we have

$$
\begin{equation*}
\sigma\left(X u_{i}\right)=D_{1}\left(X u_{i}\right) \quad \forall u_{i} \in \mathcal{K}_{1}, \tag{154}
\end{equation*}
$$

where $D_{1}$ is the $n \times n$ diagonal matrix with only the first diagonal element equal to 1 and the rest 0 . Similarly, there exist a set of $u^{\prime} s$ which result in $\sigma(X u)$ having certain components to be non-zero and the rest zero. For each such combination of zero and non-zero components, we will have a corresponding set of $u_{i}^{\prime} s$ and a corresponding $n \times n$ Diagonal matrix $D_{i}$. We define the possible set of such diagonal matrices possible for a given matrix X as

$$
\begin{equation*}
D_{X}=\left\{D=\operatorname{diag}(\mathbf{1}(X u \geq 0)): u \in \mathbb{R}^{d}, D \in \mathbb{R}^{n \times n}\right\} \tag{155}
\end{equation*}
$$

where $\operatorname{diag}(\mathbf{1}(X u \geq 0))$ represents a matrix given by

$$
D_{k, j}= \begin{cases}\mathbf{1}\left(x_{j}^{T} u\right), & \text { for } k=j  \tag{156}\\ 0 & \text { for } k \neq j\end{cases}
$$

where $\mathbf{1}(x)=1$ if $x>0$ and $\mathbf{1}(x)=0$ if $x \leq 0$. Corresponding to each such matrix $D_{i}$, there exists a set of $u_{i}$ given by

$$
\begin{equation*}
\mathcal{K}_{i}=\left\{u \in \mathbb{R}^{d}: \sigma\left(X u_{i}\right)=D_{i} X u_{i}, D_{i} \in D_{X}\right\} \tag{157}
\end{equation*}
$$

where $I$ is the $n \times n$ identity matrix. The number of these matrices $D_{i}$ is upper bounded by $2^{n}$. From Wang et al. (2021b) the upper bound is $\mathcal{O}\left(r\left(\frac{n}{r}\right)^{r}\right)$ where $r=\operatorname{rank}(X)$. Also, note that the sets $\mathcal{K}_{i}$ form a partition of the space $\mathbb{R}^{d \times 1}$. Using these definitions, we define the equivalent convex problem to the one in Equation equation 147 as

$$
\begin{equation*}
\mathcal{L}_{\beta}(v, w):=\left(\left\|\sum_{D_{i} \in D_{X}} D_{i}\left(X\left(v_{i}-w_{i}\right)\right)-y\right\|_{2}^{2}\right), \tag{158}
\end{equation*}
$$

where $v=\left\{v_{i}\right\}_{i \in 1, \cdots,\left|D_{X}\right|}, w=\left\{w_{i}\right\}_{i \in 1, \cdots,\left|D_{X}\right|}, v_{i}, w_{i} \in \mathcal{K}_{i}$, note that by definition, for any fixed $i \in\left\{1, \cdots,\left|D_{X}\right|\right\}$ at-least one of $v_{i}$ or $w_{i}$ are zero. If $v^{*}, w^{*}$ are the optimal solutions to Equation equation 158, the number of neurons $m$ of the original problem in Equation equation 147 should be greater than the number of elements of $v^{*}, w^{*}$, which have at-least one of $v_{i}^{*}$ or $w_{i}^{*}$ non-zero. We denote this value as $m_{X, y}^{*}$, with the subscript $X$ denoting that this quantity depends upon the data matrix $X$ and response $y$.
We convert $v^{*}, w^{*}$ to optimal values of Equation equation 147, denoted by $\theta^{*}=\left(U^{*}, \alpha^{*}\right)$, using a function $\psi: \mathbb{R}^{d} \times \mathbb{R}^{d} \rightarrow \mathbb{R}^{d} \times \mathbb{R}$ defined as follows

$$
\psi\left(v_{i}, w_{i}\right)=\left\{\begin{array}{lr}
\left(v_{i}, 1\right), & \text { if } w_{i}=0  \tag{159}\\
\left(w_{i},-1\right), & \text { if } v_{i}=0 \\
(0,0), & \text { if } v_{i}=w_{i}=0
\end{array}\right.
$$

where according to Pilanci \& Ergen (2020) we have $\left(u_{i}^{*}, \alpha_{i}^{*}\right)=\psi\left(v_{i}^{*}, w_{i}^{*}\right)$, for all $i \in$ $\left\{1, \cdots,\left|D_{X}\right|\right\}$ where $u_{i}^{*}, \alpha_{i}^{*}$ are the elements of $\theta^{*}$. Note that restriction of $\alpha_{i}$ to $\{1,-1,0\}$ is shown to be valid in Mishkin et al. (2022). For $i \in\left\{\left|D_{X}\right|+1, \cdots, m\right\}$ we set $\left(u_{i}^{*}, \alpha_{i}^{*}\right)=(0,0)$.
Since $D_{X}$ is hard to obtain computationally unless $X$ is of low rank, we can construct a subset $\tilde{D} \in D_{X}$ and perform the optimization in Equation equation 158 by replacing $D_{X}$ with $\tilde{D}$ to get

$$
\begin{equation*}
\mathcal{L}_{\beta}(v, w):=\left(\left\|\sum_{D_{i} \in \tilde{D}} D_{i}\left(X\left(v_{i}-w_{i}\right)\right)-y\right\|_{2}^{2}\right) \tag{160}
\end{equation*}
$$

where $v=\left\{v_{i}\right\}_{i \in 1, \cdots,|\tilde{D}|}, w=\left\{w_{i}\right\}_{i \in 1, \cdots,|\tilde{D}|}, v_{i}, w_{i} \in \mathcal{K}_{i}$, by definition, for any fixed $i \in$ $\{1, \cdots,|\tilde{D}|\}$ at-least one of $v_{i}$ or $w_{i}$ are zero.
The required condition for $\tilde{D}$ to be a sufficient replacement for $D_{X}$ is as follows. Suppose $(v, w)=$ $\left(\bar{v}_{i}, \bar{w}_{i}\right)_{i \in(1, \cdots,|\tilde{D}|)}$ denote the optimal solutions of Equation equation 160 . Then we require

$$
\begin{equation*}
m \geq \sum_{D_{i} \in \tilde{D}}\left|\left\{\bar{v}_{i}: \bar{v}_{i} \neq 0\right\} \cup\left\{\bar{w}_{i}: \bar{w}_{i} \neq 0\right\}\right| . \tag{161}
\end{equation*}
$$

Or, the number of neurons in the neural network are greater than the number of indices $i$ for which at-least one of $v_{i}^{*}$ or $w_{i}^{*}$ is non-zero. Further,

$$
\begin{equation*}
\operatorname{diag}\left(X u_{i}^{*} \geq 0: i \in[m]\right) \in \tilde{D} \tag{162}
\end{equation*}
$$

In other words, the diagonal matrices induced by the optimal $u_{i}^{*}$ 's of Equation equation 147 must be included in our sample of diagonal matrices. This is proved in Theorem 2.1 of Mishkin et al. (2022).
A computationally efficient method for obtaining $\tilde{D}$ and obtaining the optimal values of the Equation equation 147, is laid out in Mishkin et al. (2022). In this method we first get our sample of diagonal matrices $\tilde{D}$ by first sampling a fixed number of vectors from a $d$ dimensional standard multivariate
distribution, multiplying the vectors with the data matrix $X$ and then forming the diagonal matrices based of which co-ordinates are positive. Then we solve an optimization similar to the one in Equation equation 158 , without the constraints, that its parameters belong to sets of the form $\mathcal{K}_{i}$ as follows.

$$
\begin{equation*}
\mathcal{L}_{\beta}^{\prime}(p):=\left(\left\|\sum_{D_{i} \in \tilde{D}} D_{i}\left(X p_{i}\right)-y\right\|_{2}^{2}\right), \tag{163}
\end{equation*}
$$

where $p \in \mathbb{R}^{d \times|\tilde{D}|}$. In order to satisfy the constraints of the form given in Equation equation 158, this step is followed by a cone decomposition step. This is implemented through a function $\left\{\psi_{i}^{\prime}\right\}_{i \in\{1, \cdots,|\tilde{D}|\}}$. Let $p^{*}=\left\{p_{i}^{*}\right\}_{i \in\{1, \cdots,|\tilde{D}|\}}$ be the optimal solution of Equation equation 163. For each $i$ we define a function $\psi_{i}^{\prime}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d} \times \mathbb{R}^{d}$ as

$$
\begin{align*}
\psi_{i}^{\prime}\left(p_{i}\right) & =\left(v_{i}, w_{i}\right)  \tag{164}\\
\text { such that } p & =v_{i}-w_{i}, \text { and } v_{i}, w_{i} \in \mathcal{K}_{i}
\end{align*}
$$

Then we obtain $\psi\left(p_{i}^{*}\right)=\left(\bar{v}_{i}, \bar{w}_{i}\right)$. As before, at-least one of $v_{i}, w_{i}$ is 0 . Note that in practice we do not know if the conditions in Equation equation 161 and equation 162 are satisfied for a given sampled $\tilde{D}$. We express this as follows. If $\tilde{D}$ was the full set of Diagonal matrices then we would have $\left(\bar{v}_{i}, \bar{w}_{i}\right)=v_{i}^{*}, w_{i}^{*}$ and $\psi\left(\bar{v}_{i}, \bar{w}_{i}\right)=\left(u_{i}^{*}, \alpha_{i}^{*}\right)$ for all $i \in\left(1, \cdots,\left|D_{X}\right|\right)$. However, since that is not the case and $\tilde{D} \in D_{X}$, this means that $\left\{\psi\left(\bar{v}_{i}, \bar{w}_{i}\right)\right\}_{i \in(1, \cdots,|\tilde{D}|)}$ is an optimal solution of a nonconvex optimization different from the one in Equation equation 147. We denote this non-convex optimization as $\mathcal{L}_{|\tilde{D}|}(\theta)$ defined as

$$
\begin{equation*}
\mathcal{L}_{|\tilde{D}|}(\theta)=\left\|\sum_{i=1}^{m^{\prime}} \sigma\left(X u_{i}\right) \alpha_{i}-y\right\|_{2}^{2} \tag{165}
\end{equation*}
$$

where $m^{\prime}=|\tilde{D}|$ or the size of the sampled diagonal matrix set. In order to quantify the error incurred due to taking a subset of $D_{X}$, we assume that the expectation of the absolute value of the difference between the neural networks corresponding to the optimal solutions of the non-convex optimizations given in Equations equation 165 and equation 147 is upper bounded by a constant depending on the size of $\tilde{D}$. The formal assumption and its justification is given in Assumption 5.

## F. 1 Proposed Natural Actor Critic Algorithm with 2-Layer Critic Parametrization (NAC2L)

We summarize the proposed approach in Algorithm 2. Algorithm 2 has an outer for loop with two inner for loops. At a fixed iteration $k$ of the outer for loop and iteration $j$ of the first inner for loop, we obtain a sequence of state action pairs and the corresponding state and reward by following the estimate of the policy at the start of the iteration. In order to perform the critic update, the state action pairs and the corresponding target $Q$ values are stored in matrix form and passed to Algorithm 3, as the input and output values respectively to solve the following optimization problem.

$$
\begin{equation*}
\underset{\theta \in \Theta}{\arg \min } \frac{1}{n} \sum_{i=1}^{n}\left(Q_{\theta}\left(s_{i}, a_{i}\right)-r\left(s_{i}, a_{i}\right)-\gamma Q_{k, j-1}\left(s_{i+1}, a_{i+1}\right)\right)^{2}, \tag{167}
\end{equation*}
$$

where $Q_{k, j-1}$ is the estimate of the $Q$ function at the $k^{t h}$ iteration of the outer for loop and the $(j-1)^{t h}$ iteration of the first inner for loop of Algorithm 2. $Q_{\theta}$ is a neural network defined as in equation 149 and $n$ is the number of state action pairs sampled at the $k^{t h}$ iteration of the outer for loop and the $j^{\text {th }}$ iteration of the first inner for loop of Algorithm 2. This is done at each iteration of the first inner for loop to perform what is known as a Fitted Q-iteration step to obtain the estimate of the critic.
Algorithm 3 first samples a set of diagonal matrices denoted by $\tilde{D}$ in line 2 of Algorithm 3. The elements of $\tilde{D}$ act as the diagonal matrix replacement of the ReLU function. Algorithm 3 then solves an optimization of the form given in Equation equation 167 by converting it to an optimization of the form equation 152. This convex optimization is solved in Algorithm 3 using the projected

```
Algorithm 2 Natural Actor Critic with 2-Layer Critic Parametrization (NAC2L)
Input: Input: \(\mathcal{S}, \mathcal{A}, \gamma\), Time Horizon \(\mathrm{K} \in \mathcal{Z}\), Updates per time step \(\mathrm{J} \in \mathcal{Z}\),starting state sam-
pling distribution \(\nu\), Actor step sizes \(\beta_{i, k}, \forall k \in\{1, \cdots, K\}, i \in\{1, \cdots, n . J\}\), Critic step size
\(\alpha\), policy gradient step size \(\eta\), Number of convex optimization steps \(T_{k, j}, k \in\{1, \cdots, K\}, j \in\)
\(\{1, \cdots, J\}\),
    Initialize: \(\lambda_{0}=\{0\}^{d}\)
    for \(k \in\{1, \cdots, K\}\) do
        \(X_{k}=\varnothing\)
        for \(j \in\{1, \cdots, J\}\) do
            Take \(n\) state action pairs sampled from \(\nu\) as the starting state distribution and then following
            policy \(\pi_{\lambda_{k}}\).
            Set \(y_{i}=r_{i}+\gamma Q_{k, j-1}\left(s_{i+1}, a_{i+1}\right)\), where \(i \in\{1, \cdots, n\}\)
            Set \(X_{j}, Y_{j}\) as the matrix of the sampled state action pairs and vector of estimated \(Q\) values
            respectively
            \(X_{k}=X_{k} \cup X_{j}\)
            Call Algorithm 3 with input ( \(X=X_{j}, y=Y_{j}, T=T_{k, j}\) ) and return parameter \(\theta\)
            \(Q_{k, j}=Q_{\theta}\)
        end for
        \(w_{0}=0^{d}\)
        for \(i \in\left\{1, \cdots,\left|X_{k}\right|\right\}\) do
            \(A_{k, J}\left(s_{i}, a_{i}\right)=Q_{k, J}\left(s_{i}, a_{i}\right)-\sum_{a \in \mathcal{A}} \pi_{\lambda_{k}}\left(a \mid s_{i}\right) Q_{k, J}\left(s_{i}, a\right)\)
            \(w_{i}=w_{i}-\beta_{i}\left(w_{i} \cdot \nabla_{\lambda} \log \pi_{\lambda_{k}}\left(a_{i} \mid s_{i}\right)-A_{k, J}\left(s_{i}, a_{i}\right)\right) \nabla_{\lambda} \log \pi_{\lambda_{k}}\left(a_{i} \mid s_{i}\right)\)
        end for
        Update \(\lambda_{k+1}=\lambda_{k}+\eta\left(\frac{1}{1-\gamma}\right) w_{\left|X_{k}\right|}\)
    end for
    Output: \(\pi_{\lambda_{K+1}}\)
```

```
Algorithm 3 Neural Network Parameter Estimation
    Input: data \((X, y, T)\)
    Sample: \(\tilde{D}=\operatorname{diag}\left(1\left(X g_{i}>0\right)\right): g_{i} \sim \mathcal{N}(0, I), i \in[|\tilde{D}|]\)
    Initialize \(y^{1}=0, u^{1}=0\)
    Initialize \(g(u)=\left\|\sum_{D_{i} \in \tilde{D}} D_{i} X u_{i}-y\right\|_{2}^{2}\)
    for \(i \in\{0, \cdots, T\}\) do
        \(u^{i+1}=y_{i}-\alpha_{i} \nabla g\left(y_{i}\right)\)
        \(y^{i+1}=\arg \min _{y:|y|_{1} \leq \frac{R_{\max }}{1-\gamma}}\left\|u_{i+1}-y\right\|_{2}^{2}\)
    end for
    Set \(u^{T+1}=u^{*}\)
    Solve Cone Decomposition:
    \(\left.\bar{v}, \bar{w} \in u_{i}^{*}=v_{i}-w_{i}, i \in[d]\right\}\) such that \(v_{i}, w_{i} \in \mathcal{K}_{i}\) and at-least one \(v_{i}, w_{i}\) is zero.
10: Construct \(\left(\theta=\left\{u_{i}, \alpha_{i}\right\}\right)\) using the transformation
```

$$
\psi\left(v_{i}, w_{i}\right)=\left\{\begin{array}{lr}
\left(v_{i}, 1\right), & \text { if } w_{i}=0  \tag{166}\\
\left(w_{i},-1\right), & \text { if } v_{i}=0 \\
(0,0), & \text { if } v_{i}=w_{i}=0
\end{array}\right.
$$

for all $i \in\{1, \cdots, m\}$
Return $\theta$
gradient descent algorithm. After obtaining the optima for this convex program, denoted by $u^{*}=$ $\left\{u_{i}^{*}\right\}_{i \in\{1, \cdots,|\tilde{D}|\}}$, in line 10 , we transform them into an estimate of the solutions for the optimization given in equation 167, which are then passed back to Algorithm 2. The procedure is described in detail along with the relevant definitions in Appendix F.

The estimate of $w_{k}^{*}$ is obtained in the second inner for loop of Algorithm equation 2 where a gradient descent is performed for the loss function of the form given in Equation equation 9 using the state
action pairs sampled in the first inner for loop. Note that we do not have access to the true $Q$ function that is required for the critic update. Thus we use the estimate of the $Q$ function obtained at the end of the first inner for loop. After obtaining our estimate of the minimizer of Equation equation 9, we update the policy parameter using the stochastic gradient update step. Here the state action pairs used are the same we sampled in the first inner for loop.

Now since we do not know the exact convex reformulation but instead use an approximation, we have the following additional assumption.
Assumption 5. Let $\theta^{*} \triangleq \arg \min _{\theta \in \Theta} \mathcal{L}(\theta)$, where $\mathcal{L}(\theta)$ is defined in equation 147 and we denote $Q_{\theta^{*}}(\cdot)$ as $Q_{\theta}(\cdot)$ as defined in equation 149 for $\theta=\theta^{*}$. Also, let $\theta_{\tilde{D}}^{*} \triangleq \arg \min _{\theta \in \Theta} \mathcal{L}_{|\tilde{D}|}(\theta)$, where $\mathcal{L}_{\tilde{D}}(\theta)$ is the loss function $\mathcal{L}(\theta)$ with the set of diagonal matrices $D$ replaced by $\tilde{D} \in D$. Further, we denote $Q_{\theta_{|\tilde{D}|}^{*}}(\cdot)$ as $Q_{\theta}(\cdot)$ as defined in equation 149 for $\theta=\theta_{|\tilde{D}|}^{*}$. Then we assume

$$
\begin{equation*}
\mathbb{E}_{s, a}\left(\left|Q_{\theta^{*}}-Q_{\theta_{|\tilde{D}|}^{*}}\right|\right)_{\nu} \leq \epsilon_{|\tilde{D}|} \tag{168}
\end{equation*}
$$

for any $(s, a) \sim \zeta_{\pi^{\lambda} k}^{\nu}$.
Thus, $\epsilon_{|\tilde{D}|}$ is a measure of the error incurred due to taking a sample of diagonal matrices $\tilde{D}$ and not the full set $D_{X}$. In practice, setting $|\tilde{D}|$ to be the same order of magnitude as $d$ (dimension of the data) gives us a sufficient number of diagonal matrices to get a reformulation of the non convex optimization problem which performs comparably or better than existing gradient descent algorithms, therefore $\epsilon_{|\tilde{D}|}$ is only included for theoretical completeness and will be negligible in practice. This has been practically demonstrated in Mishkin et al. (2022); Bartan \& Pilanci (2022); Sahiner et al. (2022).

Before proceeding with the proof further, we would like to prove two supplementary lemmas
Lemma 8. Consider an optimization of the form given in Equation equation 160 denoted by $\mathcal{L}_{|\tilde{D}|}$ and it's convex equivalent denoted by $\mathcal{L}_{0}$. Then the value of these two loss functions evaluated at $(v, w)=\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}$ and $\theta=\psi\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}$ respectively are equal and thus we have

$$
\begin{equation*}
\mathcal{L}_{|\tilde{D}|}\left(\psi\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}\right)=\mathcal{L}_{0}\left(\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}\right) \tag{169}
\end{equation*}
$$

Proof. Consider the loss functions in Equations equation 158, equation 163 are as follows

$$
\begin{align*}
\mathcal{L}_{0}\left(\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}\right) & =\left\|\sum_{D_{i} \in \tilde{D}} D_{i}\left(X\left(v_{i}-w_{i}\right)\right)-y\right\|_{2}^{2}  \tag{170}\\
\mathcal{L}_{|\tilde{D}|}\left(\psi\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}\right) & =\left\|\sum_{i=1}^{|\tilde{D}|} \sigma\left(X \psi\left(v_{i}, w_{i}\right)_{1}\right) \psi\left(v_{i}, w_{i}\right)_{2}-y\right\|_{2}^{2}, \tag{171}
\end{align*}
$$

where $\psi\left(v_{i}, w_{i}\right)_{1}, \psi\left(v_{i}, w_{i}\right)_{2}$ represent the first and second coordinates of $\psi\left(v_{i}, w_{i}\right)$ respectively.
For any fixed $i \in\{1, \cdots,|\tilde{D}|\}$ consider the two terms

$$
\begin{array}{r}
D_{i}\left(X\left(v_{i}-w_{i}\right)\right) \\
\sigma\left(X \psi\left(v_{i}, w_{i}\right)_{1}\right) \psi\left(v_{i}, w_{i}\right)_{2} \tag{173}
\end{array}
$$

For a fixed $i$ either $v_{i}$ or $w_{i}$ is zero. In case both are zero, both of the terms in Equations equation 172 and equation 173 are zero as $\psi(0,0)=(0,0)$. Assume that for a given $i w_{i}=0$. Then we have $\psi\left(v_{i}, w_{i}\right)=\left(v_{i}, 1\right)$. Then equations equation 172, equation 173 are.

$$
\begin{gather*}
D_{i}\left(X\left(v_{i}\right)\right.  \tag{174}\\
\sigma\left(X\left(v_{i}\right)\right) \tag{175}
\end{gather*}
$$

But by definition of $v_{i}$ we have $D_{i}\left(X\left(v_{i}\right)=\sigma\left(X\left(v_{i}\right)\right)\right.$, therefore Equations equation 174, equation 175 are equal. Alternatively if for a given $i v_{i}=0$, then $\psi\left(v_{i}, w_{i}\right)=\left(w_{i},-1\right)$, then the terms in equation 172, equation 173 become.

$$
\begin{align*}
& -D_{i}\left(X\left(w_{i}\right)\right.  \tag{176}\\
& -\sigma\left(X\left(w_{i}\right)\right) \tag{177}
\end{align*}
$$

By definition of $w_{i}$ we have $D_{i}\left(X\left(w_{i}\right)=\sigma\left(X\left(w_{i}\right)\right)\right.$, then the terms in equation 176, equation 176 are equal. Since this is true for all $i$, we have

$$
\begin{equation*}
\mathcal{L}_{|\tilde{D}|}\left(\psi\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}\right)=\mathcal{L}_{0}\left(\left(v_{i}, w_{i}\right)_{i \in\{1, \cdots,|\tilde{D}|\}}\right) \tag{178}
\end{equation*}
$$

Lemma 9. The function $Q_{\theta}(x)$ defined in equation equation 149 is Lipschitz continuous in $\theta$, where $\theta$ is considered a vector in $\mathbb{R}^{(d+1) m}$ with the assumption that the set of all possible $\theta$ belong to the set $\mathcal{B}=\left\{\theta:\left|\theta^{*}-\theta\right|_{1}<1\right\}$, where $\theta^{*}$ is some fixed value.

Proof. First we show that for all $\theta_{1}=\left\{u_{i}, \alpha_{i}\right\}, \theta_{2}=\left\{u_{i}^{\prime}, \alpha_{i}^{\prime}\right\} \in \mathcal{B}$ we have $\alpha_{i}=\alpha_{i}^{\prime}$ for all $i \in(1, \cdots, m)$
Note that

$$
\begin{equation*}
\left|\theta_{1}-\theta_{2}\right|_{1}=\sum_{i=1}^{m}\left|u_{i}-u_{i}^{\prime}\right|_{1}+\sum_{i=1}^{m}\left|\alpha_{i}-\alpha_{i}^{\prime}\right| \tag{179}
\end{equation*}
$$

where $\left|u_{i}-u_{i}^{\prime}\right|_{1}=\sum_{j=1}^{d}\left|u_{i_{j}}-u_{i_{j}}^{\prime}\right|$ with $u_{i_{j}}, u_{i_{j}}^{\prime}$ denote the $j^{t h}$ component of $u_{i}, u_{i}^{\prime}$ respectively. By construction $\alpha_{i}, \alpha_{i}^{\prime}$ can only be $1,-1$ or 0 . Therefore if $\alpha_{i} \neq \alpha_{i}^{\prime}$ then $\left|\alpha_{i}-\alpha_{i}^{\prime}\right|=2$ if both non zero or $\left|\alpha_{i}-\alpha_{i}^{\prime}\right|=1$ if one is zero. Therefore $\left|\theta_{1}-\theta_{2}\right|_{1} \geq 1$. Which leads to a contradiction.

Therefore $\alpha_{i}=\alpha_{i}^{\prime}$ for all $i$ and we also have

$$
\begin{equation*}
\left|\theta_{1}-\theta_{2}\right|_{1}=\sum_{i=1}^{m}\left|u_{i}-u_{i}^{\prime}\right|_{1} \tag{180}
\end{equation*}
$$

$Q_{\theta}(x)$ is defined as

$$
\begin{equation*}
Q_{\theta}(x)=\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i} \tag{181}
\end{equation*}
$$

From Proposition 1 in Scaman \& Virmaux (2018) the function $Q_{\theta}(x)$ is Lipschitz continuous in $x$, therefore there exist $l>0$ such that

$$
\begin{align*}
\left|Q_{\theta}(x)-Q_{\theta}(y)\right| & \leq l|x-y|_{1}  \tag{182}\\
\left|\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i}-\sum_{i=1}^{m} \sigma^{\prime}\left(y^{T} u_{i}\right) \alpha_{i}\right| & \leq l|x-y|_{1} \tag{183}
\end{align*}
$$

If we consider a single neuron of $Q_{\theta}$, for example $i=1$, we have $l_{1}>0$ such that

$$
\begin{equation*}
\left|\sigma^{\prime}\left(x^{T} u_{1}\right) \alpha_{i}-\sigma^{\prime}\left(y^{T} u_{1}\right) \alpha_{i}\right| \leq l_{1}|x-y|_{1} \tag{184}
\end{equation*}
$$

Now consider Equation equation 184, but instead of considering the left hand side a a function of $x, y$ consider it a function of $u$ where we consider the difference between $\sigma^{\prime}\left(x^{T} u\right) \alpha_{i}$ evaluated at $u_{1}$ and $u_{1}^{\prime}$ such that

$$
\begin{equation*}
\left|\sigma^{\prime}\left(x^{T} u_{1}\right) \alpha_{i}-\sigma^{\prime}\left(x^{T} u_{1}^{\prime}\right) \alpha_{i}\right| \leq l_{1}^{x}\left|u_{1}-u_{1}^{\prime}\right|_{1} \tag{185}
\end{equation*}
$$

for some $l_{1}^{x}>0$.
Similarly, for all other $i$ if we change $u_{i}$ to $u_{i}^{\prime}$ to be unchanged we have

$$
\begin{equation*}
\left|\sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i}-\sigma^{\prime}\left(x^{T} u_{i}^{\prime}\right) \alpha_{i}\right| \leq l_{i}^{x}\left|u_{i}-u_{i}^{\prime}\right|_{1} \tag{186}
\end{equation*}
$$

for all $x$ if both $\theta_{1}, \theta_{2} \in \mathcal{B}$.
Therefore we obtain

$$
\begin{align*}
\left|\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i}-\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}^{\prime}\right) \alpha_{i}\right| & \leq \sum_{i=1}^{m}\left|\sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i}-\left(x^{T} u_{i}^{\prime}\right) \alpha_{i}\right|  \tag{187}\\
& \leq \sum_{i=1}^{m} l_{i}^{x}\left|u_{i}-u_{i}^{\prime}\right|_{1}  \tag{188}\\
& \leq\left(\sup _{i} l_{i}^{x}\right) \sum_{i=1}^{m}\left|u_{i}-u_{i}^{\prime}\right|_{1}  \tag{189}\\
& \leq\left(\sup _{i} l_{i}^{x}\right)\left|\theta_{1}-\theta_{2}\right| \tag{190}
\end{align*}
$$

This result for a fixed $x$. If we take the supremum over $x$ on both sides we get

$$
\begin{equation*}
\sup _{x}\left|\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i}-\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}^{\prime}\right) \alpha_{i}\right| \leq\left(\sup _{i, x} l_{i}^{x}\right)\left|\theta_{1}-\theta_{2}\right| \tag{191}
\end{equation*}
$$

Denoting $\left(\sup _{i, x} l_{i}^{x}\right)=l$, we get

$$
\begin{align*}
\left|\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}\right) \alpha_{i}-\sum_{i=1}^{m} \sigma^{\prime}\left(x^{T} u_{i}^{\prime}\right) \alpha_{i}\right| \leq & l\left|\theta_{1}-\theta_{2}\right|_{1}  \tag{192}\\
& \forall x \in \mathbb{R}^{d} \tag{193}
\end{align*}
$$

Now since the critic step becomes a convex optimization Lemma 6 now becomes
Lemma 10. For a given iteration $k$ of Algorithm 2 and iteration $j$ of the first for loop of Algorithm 2, let the number of steps of the projected gradient descent performed by Algorithm 3, denoted by $T_{k, j}$, and the gradient descent step size $\alpha_{k, j}$ satisfy

$$
\begin{equation*}
\alpha_{k, j}=\frac{\left\|u_{k, j}^{*}\right\|_{2}}{L_{k, j} \sqrt{T_{k, j}+1}}, \tag{194}
\end{equation*}
$$

for some constants $L_{k, j}$ and $\left\|\left(u_{k}^{*}\right)\right\|_{2}$. Then the error $\epsilon_{k_{4}}$ defined in Definition 7 is upper bounded as

$$
\begin{equation*}
\mathbb{E}\left(\left|\epsilon_{k, j}^{4}\right|\right) \leq \tilde{\mathcal{O}}\left(\frac{1}{\sqrt{T_{k, j}}}\right)+\epsilon_{|\tilde{D}|} \tag{195}
\end{equation*}
$$

Proof. For a given iteration $k$ of Algorithm 1 and iteration $j$ of the first inner for loop, the optimization problem to be solved in Algorithm 3 is the following

$$
\begin{equation*}
\mathcal{L}(\theta)=\frac{1}{n} \sum_{i=1}^{n}\left(Q_{\theta}\left(s_{i}, a_{i}\right)-\left(r\left(s_{i}, a_{i}\right)+\gamma \max _{a^{\prime} \in \mathcal{A}} \gamma Q_{k, j-1}\left(s^{\prime}, a^{\prime}\right)\right)\right)^{2} \tag{196}
\end{equation*}
$$

Here, $Q_{k, j-1}$ is the estimate of the $Q$ function from the iteration $j-1$ and the state action pairs $\left(s_{i}, a_{i}\right)_{i=\{1, \cdots, n\}}$ have been sampled from a distribution over the state action pairs denoted by $\nu$. Since $\min _{\theta} \mathcal{L}(\theta)$ is a non convex optimization problem we instead solve the equivalent convex problem given by

$$
\begin{align*}
u_{k, j}^{*}= & \underset{u}{\arg \min } g_{k, j}(u)=\underset{u}{\arg \min }\left\|\sum_{D_{i} \in \tilde{D}} D_{i} X_{k, j} u_{i}-y_{k}\right\|_{2}^{2}  \tag{197}\\
& \text { subject to }|u|_{1} \leq \frac{R_{\max }}{1-\gamma} \tag{198}
\end{align*}
$$

Here, $X_{k, j} \in \mathbb{R}^{n \times d}$ is the matrix of sampled state action pairs at iteration $k, y_{k} \in \mathbb{R}^{n \times 1}$ is the vector of target values at iteration $k . \tilde{D}$ is the set of diagonal matrices obtained from line 2 of Algorithm 3 and $u \in \mathbb{R}^{|\tilde{D} d| \times 1}$ (Note that we are treating $u$ as a vector here for notational convenience instead of a matrix as was done in Section 4).
The constraint in Equation equation 198 ensures that the all the co-ordinates of the vector $\sum_{D_{i} \in \tilde{D}} D_{i} X_{k, j} u_{i}$ are upper bounded by $\frac{R_{\max }}{1-\gamma}$ (since all elements of $X_{k, j}$ are between 0 and 1 ). This ensures that the corresponding neural network represented by Equation equation 149 is also upper bounded by $\frac{R_{\max }}{1-\gamma}$. We use the a projected gradient descent to solve the constrained convex optimization problem which can be written as.

$$
\begin{equation*}
u_{k, j}^{*}=\underset{u:|u|_{1} \leq \frac{R_{\text {max }}}{1-\gamma}}{\arg \min } g_{k, j}(u)=\underset{u:|u|_{1} \leq \frac{R_{\text {max }}}{1-\gamma}}{\arg \min }\left\|\sum_{D_{i} \in \tilde{D}} D_{i} X_{k, j} u_{i}-y_{k}\right\|_{2}^{2} \tag{199}
\end{equation*}
$$

From Ang, Andersen(2017). "Continuous Optimization" [Notes]. https://angms.science/doc/CVX we have that if the step size $\alpha=\frac{\left\|u_{k, j}^{*}\right\|_{2}}{L_{k, j} \sqrt{T_{k, j}+1}}$, after $T_{k, j}$ iterations of the projected gradient descent algorithm we obtain

$$
\begin{equation*}
\left(g_{k, j}\left(u_{T_{k, j}}\right)-g_{k, j}\left(u^{*}\right)\right) \leq L_{k, j} \frac{\left\|u_{k, j}^{*}\right\|_{2}}{\sqrt{T_{k, j}+1}} \tag{200}
\end{equation*}
$$

Where $L_{k, j}$ is the lipschitz constant of $g_{k, j}(u)$ and $u_{T_{k, j}}$ is the parameter estimate at step $T_{k, j}$.
Therefore if the number of iteration of the projected gradient descent algorithm $T_{k, j}$ and the step-size $\alpha$ satisfy

$$
\begin{align*}
T_{k, j} & \geq L_{k, j}^{2}\left\|u_{k, j}^{*}\right\|_{2}^{2} \epsilon^{-2}-1  \tag{201}\\
\alpha & =\frac{\left\|u_{k, j}^{*}\right\|_{2}}{L_{k, j} \sqrt{T_{k, j}+1}} \tag{202}
\end{align*}
$$

we have

$$
\begin{equation*}
\left(g_{k, j}\left(u_{T_{k, j}}\right)-g_{k, j}\left(u^{*}\right)\right) \leq \epsilon \tag{203}
\end{equation*}
$$

Let $\left(v_{i}^{*}, w_{i}^{*}\right)_{i \in(1, \cdots,|\tilde{D}|)},\left(v_{i}^{T_{k, j}}, w_{i}^{T_{k, j}}\right)_{i \in(1, \cdots,|\tilde{D}|)}$ be defined as

$$
\begin{array}{r}
\left(v_{i}^{*}, w_{i}^{*}\right)_{i \in(1, \cdots,|\tilde{D}|)}=\psi_{i}^{\prime}\left(u_{i}^{*}\right)_{i \in(1, \cdots,|\tilde{D}|)} \\
\left(v_{i}^{T_{k, j}}, w_{i}^{T_{k, j}}\right)_{i \in(1, \cdots,|\tilde{D}|)}=\psi_{i}^{\prime}\left(u_{i}^{T_{k, j}}\right)_{i \in(1, \cdots,|\tilde{D}|)} \tag{205}
\end{array}
$$

where $\psi^{\prime}$ is defined in Equation equation 164.
Further, we define $\theta_{|\tilde{D}|}^{*}$ and $\theta^{T_{k, j}}$ as

$$
\begin{array}{r}
\theta_{|\tilde{D}|}^{*}=\psi\left(v_{i}^{*}, w_{i}^{*}\right)_{i \in(1, \cdots,|\tilde{D}|)} \\
\theta^{T_{k, j}}=\psi\left(v_{i}^{T_{k, j}}, w_{i}^{T_{k, j}}\right)_{i \in(1, \cdots,|\tilde{D}|)} \tag{207}
\end{array}
$$

where $\psi$ is defined in Equation equation 159, $\theta_{|\tilde{D}|}^{*}=\arg \min _{\theta} \mathcal{L}_{|\tilde{D}|}(\theta)$ for $\mathcal{L}_{|\tilde{D}|}(\theta)$ defined in Appendix F.
Since $\left(g\left(u_{T_{k, j}}\right)-g\left(u^{*}\right)\right) \leq \epsilon$, then by Lemma 8 , we have

$$
\begin{equation*}
\mathcal{L}_{|\tilde{D}|}\left(\theta^{T_{k, j, j}}\right)-\mathcal{L}_{|\tilde{D}|}\left(\theta_{|\tilde{D}|}^{*}\right) \leq \epsilon \tag{208}
\end{equation*}
$$

Note that $\mathcal{L}_{|\tilde{D}|}\left(\theta^{T_{k, j}}\right)-\mathcal{L}_{|\tilde{D}|}\left(\theta_{|\tilde{D}|}^{*}\right)$ is a constant value. Thus we can always find constant $C_{k, j}^{\prime}$ such that

$$
\begin{gather*}
C_{k}^{\prime}\left|\theta^{T_{k, j}}-\theta_{|\tilde{D}|}^{*}\right|_{1} \leq \mathcal{L}_{|\tilde{D}|}\left(\theta^{T_{k, j}}\right)-\mathcal{L}_{|\tilde{D}|}\left(\theta_{|\tilde{D}|}^{*}\right)  \tag{209}\\
\left|\theta^{T_{k, j}}-\theta_{|\tilde{D}|}^{*}\right|_{1} \leq \frac{\mathcal{L}\left(\theta^{T_{k, j}}\right)-\mathcal{L}\left(\theta^{*}\right)}{C_{k}^{\prime}} \tag{210}
\end{gather*}
$$

Therefore if we have

$$
\begin{align*}
T_{k, j} & \geq L_{k, j}^{2}\left\|u_{k, j}^{*}\right\|_{2}^{2} \epsilon^{-2}-1  \tag{211}\\
\alpha_{k, j} & =\frac{\left\|u_{k}^{*}\right\|_{2}}{L_{k, j} \sqrt{T_{k, j}+1}} \tag{212}
\end{align*}
$$

then we have

$$
\begin{equation*}
\left|\theta^{T_{k, j}}-\theta^{*}\right|_{1} \leq \frac{\epsilon}{C_{k}^{\prime}} \tag{213}
\end{equation*}
$$

which according to Equation equation 210 implies that

$$
\begin{equation*}
C_{k}^{\prime}\left|\theta^{T_{k, j}}-\theta_{|\tilde{D}|}^{*}\right|_{1} \leq \mathcal{L}_{|\tilde{D}|}\left(\theta^{T_{k, j}}\right)-\mathcal{L}_{|\tilde{D}|}\left(\theta_{|\tilde{D}|}^{*}\right) \leq \epsilon \tag{214}
\end{equation*}
$$

Dividing Equation equation 214 by $C_{k}^{\prime}$ we get

$$
\begin{equation*}
\left|\theta^{T_{k, j}}-\theta_{|\tilde{D}|}^{*}\right|_{1} \leq \frac{\mathcal{L}_{|\tilde{D}|}\left(\theta^{T_{k, j}}\right)-\mathcal{L}_{|\tilde{D}|}\left(\theta_{|\tilde{D}|}^{*}\right)}{C_{k}^{\prime}} \leq \frac{\epsilon}{C_{k}^{\prime}} \tag{215}
\end{equation*}
$$

Which implies

$$
\begin{equation*}
\left|\theta^{T_{k, j}}-\theta_{|\tilde{D}|}^{*}\right|_{1} \leq \frac{\epsilon}{C_{k}^{\prime}} \tag{216}
\end{equation*}
$$

Assuming $\epsilon$ is small enough such that $\frac{\epsilon}{C_{k}^{\prime}}<1$ from lemma 9, this implies that there exists an $L_{k, j}>0$ such that

$$
\begin{align*}
\left|Q_{\theta^{T_{k}, j}}(s, a)-Q_{\theta_{|\tilde{D}|}^{*}}(s, a)\right| & \leq L_{k, j}\left|\theta^{T_{k, j}}-\theta_{|\tilde{D}|}^{*}\right|_{1}  \tag{217}\\
& \leq \frac{L_{k, j} \epsilon}{C_{k}^{\prime}} \tag{218}
\end{align*}
$$

for all $(s, a) \in \mathcal{S} \times \mathcal{A}$. Equation equation 218 implies that if

$$
\begin{align*}
T_{k, j} & \geq L_{k, j}^{2}\left\|u_{k, j}^{*}\right\|_{2}^{2} \epsilon^{-2}-1  \tag{219}\\
\alpha_{k, j} & =\frac{\left\|u_{k}^{*}\right\|_{2}}{L_{k, j} \sqrt{T_{k, j}+1}} \tag{220}
\end{align*}
$$

then we have

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{\theta^{T_{k}, j}}(s, a)-Q_{\theta_{|\tilde{D}|}^{*}}(s, a)\right|\right) \leq \frac{L_{k, j} \epsilon}{C_{k}^{\prime}} \tag{221}
\end{equation*}
$$

By definition in section $\mathrm{C} Q_{k, j}$ is our estimate of the $Q$ function at the $k^{\text {th }}$ iteration of Algorithm 1 and thus we have $Q_{\theta^{T_{k, j}}}=Q_{k, j}$ which implies that

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{k, j}(s, a)-Q_{\theta_{\tilde{D}}^{*}}(s, a)\right|\right) \leq \frac{L_{k, j} \epsilon}{C_{k}^{\prime}} \tag{222}
\end{equation*}
$$

If we replace $\epsilon$ by $\frac{C_{k, j}^{\prime} \epsilon}{L_{k, j}}$ in Equation equation 221, we get that if

$$
\begin{align*}
T_{k, j} & \geq\left(\frac{C_{k, j}^{\prime} \epsilon}{L_{k, j}}\right)^{-2} L_{k, j}^{2}\left\|u_{k, j}^{*}\right\|_{2}^{2}-1  \tag{223}\\
\alpha_{k, j} & =\frac{\left\|u_{k}^{*}\right\|_{2}}{L_{k, j} \sqrt{T_{k, j}+1}} \tag{224}
\end{align*}
$$

we have

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{k, j}(s, a)-Q_{\theta_{D}^{*}}(s, a)\right|\right) \leq \epsilon \tag{225}
\end{equation*}
$$

From Assumption 5, we have that

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{\theta^{*}}(s, a)-Q_{\theta_{\tilde{D}}^{*}}(s, a)\right|\right) \leq \epsilon_{|\tilde{D}|} \tag{226}
\end{equation*}
$$

where $\theta^{*}=\arg \min _{\theta \in \Theta} \mathcal{L}(\theta)$ and by definition of $Q_{k, j}^{3}$ in Definition 6, we have that $Q_{k, j}^{3}=Q_{\theta^{*}}$. Therefore if we have

$$
\begin{align*}
T_{k, j} & \geq\left(\frac{C_{k, j}^{\prime} \epsilon}{L_{k, j}}\right)^{-2} L_{k, j}^{2}\left\|u_{k, j}^{*}\right\|_{2}^{2}-1  \tag{227}\\
\alpha_{k, j} & =\frac{\left\|u_{k}^{*}\right\|_{2}}{L_{k, j} \sqrt{T_{k, j}+1}} \tag{228}
\end{align*}
$$

we have

$$
\begin{align*}
\mathbb{E}\left(\left|Q_{k, j}(s, a)-Q_{k, j}^{3}(s, a)\right|\right)_{\nu} & \leq \mathbb{E}\left(\left|Q_{k, j}(s, a)-Q_{\theta_{\tilde{D}}^{*}}(s, a)\right|\right)+\mathbb{E}\left(\left|Q_{k, j}^{3}(s, a)-Q_{\theta_{\tilde{D}}^{*}}(s, a)\right|\right)  \tag{229}\\
& \leq \epsilon+\epsilon_{|\tilde{D}|} \tag{230}
\end{align*}
$$

This implies

$$
\begin{equation*}
\mathbb{E}\left(\left|Q_{k, j}(s, a)-Q_{k, j}^{3}(s, a)\right|\right) \leq \tilde{\mathcal{O}}\left(\frac{1}{\sqrt{T_{k, j}}}\right)+\epsilon_{|\tilde{D}|} \tag{231}
\end{equation*}
$$

Thus replacing lemma 6 with lemma 10 we obtain the following.
Theorem 2. Suppose Assumptions 1-5 hold and we have, $\alpha_{i}=\frac{\left\|u_{k, j}^{*}\right\|_{2}}{L_{k, j} \sqrt{i+1}}, \eta=\frac{1}{\sqrt{K}}$ and $\beta_{i}=$ $\frac{2}{\mu_{k}(i+1)}$, then we obtain

$$
\begin{align*}
\min _{k \leq K}\left(V^{*}(\nu)-V^{\pi_{\lambda_{K}}}(\nu)\right) \leq & \mathcal{O}\left(\frac{1}{\sqrt{K}(1-\gamma)}\right)+\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} \sum_{j=0}^{J-1} \mathcal{O}\left(\frac{\log \log (n)}{\sqrt{n}}\right)+ \\
& +\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} \sum_{j=0}^{J-1} \mathcal{O}\left(\frac{1}{\sqrt{T_{k, j}}}\right)+\frac{1}{K(1-\gamma)} \sum_{k=1}^{K} \mathcal{O}\left(\gamma^{J}\right) \\
& +\frac{1}{1-\gamma}\left(\epsilon_{\text {bias }}+\left(\sqrt{\epsilon_{\text {approx }}}\right)+\epsilon_{|\tilde{D}|}\right) \tag{232}
\end{align*}
$$

Hence, for $K=\mathcal{O}\left(\epsilon^{-2}(1-\gamma)^{-2}\right), J=\mathcal{O}\left(\log \left(\frac{1}{\epsilon}\right)\right), n=\tilde{\mathcal{O}}\left(\epsilon^{-2}(1-\gamma)^{-2}\right)$,
$T_{k, j}=\mathcal{O}\left(\epsilon^{-2}(1-\gamma)^{-2}\right)$ we have

$$
\begin{equation*}
\min _{k \leq K}\left(V^{*}(\nu)-V^{\pi_{\lambda_{K}}}(\nu)\right) \leq \epsilon+\frac{1}{1-\gamma}\left(\epsilon_{\text {bias }}+\left(\sqrt{\epsilon_{\text {approx }}}\right)+\epsilon_{|\tilde{D}|}\right) \tag{233}
\end{equation*}
$$

which implies a sample complexity of $\sum_{k=1}^{K} \sum_{j=1}^{J}(n)=\tilde{\mathcal{O}}\left(\epsilon^{-4}(1-\gamma)^{-4}\right)$.
Note that we no longer have the probability statement of 'with a probability at least $1-\delta$. However, on the flip side, we have and extra error term in the form of $\epsilon_{|\tilde{D}|}$. This error represents the error incurred due to the inability to obtain an exact convex reformulation of the 2 layer neural network.


[^0]:    ${ }^{1}$ For a measurable set $\mathcal{X}$, let $\mathcal{P}(\mathcal{X})$ denote the set of all probability measures over $\mathcal{X}$.

[^1]:    ${ }^{2}$ Algorithmic Foundations of Learning [Lecture Notes].https://www.stats.ox.ac.uk/~rebeschi/teaching/AFoL/22/

