# RESTRICTED INVERTIBILITY AND THE BANACH-MAZUR DISTANCE TO THE CUBE 

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#### Abstract

We prove a normalized version of the restricted invertibility principle obtained by Spielman-Srivastava in [12]. Applying this result, we get a new proof of the proportional Dvoretzky-Rogers factorization theorem recovering the best current estimate. As a consequence, we also recover the best known estimate for the Banach-Mazur distance to the cube: the distance of every $n$-dimensional normed space from $\ell_{\infty}^{n}$ is at most $(2 n)^{\frac{5}{6}}$. Finally, using tools from the work of Batson-Spielman-Srivastava in [2], we give a new proof for a theorem of Kashin-Tzafriri on the norm of restricted matrices.


## 1. Introduction

Given an $n \times m$ matrix $U$, viewed as an operator from $\ell_{2}^{m}$ to $\ell_{2}^{n}$, the restricted invertibility problem asks if we can extract a large number of linearly independent columns of $U$ and provide an estimate for the norm of the restricted inverse. If we write $U_{\sigma}$ for the restriction of $U$ to the columns $U e_{i}, i \in \sigma \subset\{1, \ldots, m\}$, we want to find a subset $\sigma$, of cardinality $k$ as large as possible, such that $\left\|U_{\sigma} x\right\|_{2} \geqslant c\|x\|_{2}$ for all $x \in \mathbb{R}^{\sigma}$ and to estimate the constant $c$ (which will depend on the operator $U$ ). This question was studied by Bourgain-Tzafriri (see [4]) who obtained a result for square matrices:

Given an $n \times n$ matrix $T$ (viewed as an operator on $\ell_{2}^{n}$ ) whose columns are of norm one, there exists $\sigma \subset\{1, \ldots, n\}$ with $|\sigma| \geqslant d \frac{n}{\|T\|_{2}^{2}}$ such that $\left\|T_{\sigma} x\right\|_{2} \geqslant c\|x\|_{2}$ for all $x \in \mathbb{R}^{\sigma}$, where $d, c>0$ are absolute constants.

Here and in the rest of the paper, $\|\cdot\|_{2}$ denotes the Euclidean norm when applied to a vector and the operator norm when applied to a matrix (seen as an operator on $l_{2}$ ). For a matrix $A$, $\|A\|_{\mathrm{HS}}$ denotes the Hilbert-Schmidt norm, i.e.

$$
\|A\|_{\mathrm{HS}}=\sqrt{\operatorname{Tr}\left(A \cdot A^{*}\right)}=\left(\sum_{i}\left\|C_{i}\right\|_{2}^{2}\right)^{1 / 2}
$$

where $C_{i}$ are the columns of $A$.
Since the identity operator can be decomposed in the form $I d=\sum_{j} e_{j} e_{j}^{t}$ where $\left(e_{j}\right)$ is the canonical basis of $\mathbb{R}^{n}$, the previous result states that one can find a large part of this basis (of cardinality greater than $\left.d \frac{n}{\|T\|_{2}^{2}}\right)$ on the span of which the operator $T$ is invertible and the norm of its inverse is controlled by an absolute constant.

Vershynin generalized this result for any decomposition of the identity (see [16]) and improved the estimate for the size of the subset. Using a technical iteration scheme based on the previous result of Bourgain-Tzafriri, combined with a theorem of Kashin-Tzafriri which we will discuss in the last section, he obtained the following :

Let $I d=\sum_{j \leqslant m} x_{j} x_{j}^{t}$ and let $T$ be a linear operator on $\ell_{2}^{n}$. For any $\varepsilon \in(0,1)$ one can find $\sigma \subset\{1, \ldots, m\}$ with

$$
|\sigma| \geqslant(1-\varepsilon) \frac{\|T\|_{\mathrm{HS}}^{2}}{\|T\|_{2}^{2}}
$$

such that
for all scalars $\left(a_{j}\right)$.
One can easily check that, in the case of the canonical decomposition, this is a generalization of the Bourgain-Tzafriri theorem, which was previously only proved for a fixed value of $\varepsilon$. The constant $c(\varepsilon)$ plays a crucial role in applications and we will be able to improve its value significantly (see Proposition 2.1 for the precise statement).

Back to the original restricted invertibility problem, a recent work of Spielman-Srivastava (see [12]) provides the best known estimate for the norm of the inverse matrix. Their proof uses a new deterministic method based on linear algebra, while the previous works on the subject employed probabilistic, combinatorial and functional-analytic arguments.

More precisely, Spielman-Srivastava proved the following:
Theorem 1.1. Let $x_{1}, \ldots x_{m} \in \mathbb{R}^{n}$ such that $I d=\sum_{i} x_{i} x_{i}^{t}$ and let $0<\varepsilon<1$. For every linear operator $T: \ell_{2}^{n} \rightarrow \ell_{2}^{n}$ there exists a subset $\sigma \subset\{1, \ldots, m\}$ of size $|\sigma| \geqslant\left\lfloor(1-\varepsilon)^{2} \frac{\|T\|_{\text {HS }}^{2}}{\|T\|_{2}^{2}}\right\rfloor$ for which $\left\{T x_{i}\right\}_{i \in \sigma}$ is linearly independent and

$$
\lambda_{\min }\left(\sum_{i \in \sigma}\left(T x_{i}\right)\left(T x_{i}\right)^{t}\right)>\frac{\varepsilon^{2}\|T\|_{\mathrm{HS}}^{2}}{m},
$$

where $\lambda_{\min }$ is computed on $\operatorname{span}\left\{T x_{i}\right\}_{i \in \sigma}$ or simply here $\lambda_{\min }$ denotes the smallest nonzero eigenvalue of the corresponding operator. .

One can view the previous result as an invertibility theorem for rectangular matrices. Given, as above, a decomposition of the identity and a linear operator $T$, we can associate to these an $n \times m$ matrix $U$ whose columns are the vectors $\left(T x_{j}\right)_{j \leqslant m}$. Since $I d=\sum_{j} x_{j} x_{j}^{t}$, one can easily check that

$$
U \cdot U^{t}=T \cdot T^{t}=\sum_{j}\left(T x_{j}\right) \cdot\left(T x_{j}\right)^{t}
$$

Hence,

$$
\|U\|_{\mathrm{HS}}=\|T\|_{\mathrm{HS}} \quad \text { and } \quad\|U\|_{2}=\|T\|_{2}
$$

and thus the previous result can be written in terms of the rectangular matrix $U$.
In the applications, one might need to extract multiples of the columns of the matrix. Adapting the proof of Spielman-Srivastava, we will generalize the restricted invertibility theorem for any rectangular matrix and, under some conditions, for any choice of multiples.

If $D$ is an $m \times m$ diagonal matrix with diagonal entries $\left(\alpha_{j}\right)_{j \leqslant m}$, we set $I_{D}:=\{j \leqslant m \mid$ $\left.\alpha_{j} \neq 0\right\}$ and write $D_{\sigma}^{-1}$ for the restricted inverse of $D$ i.e the diagonal matrix whose diagonal entries are the inverses of the respective entries of $D$ for indices in $\sigma$ and zero elsewhere. The main result of this paper is the following:

Theorem 1.2. Given an $n \times m$ matrix $U$ and a diagonal $m \times m$ matrix $D$ with $\left(\alpha_{j}\right)_{j \leqslant m}$ on its diagonal, with the property that $\operatorname{Ker}(D) \subset \operatorname{Ker}(U)$, then for any $\varepsilon \in(0,1)$ there exists $\sigma \subset\{1, \ldots, m\}$ with

$$
|\sigma| \geqslant(1-\varepsilon)^{2} \frac{\|U\|_{\mathrm{HS}}^{2}}{\|U\|_{2}^{2}}
$$

such that

$$
s_{\min }\left(U_{\sigma} D_{\sigma}^{-1}\right)>\frac{\varepsilon\|U\|_{\mathrm{HS}}}{\|D\|_{\mathrm{HS}}}
$$

where $s_{\min }$ denotes the smallest singular value.

Note that if we apply this fact to the matrix $U$ which we associated with a linear operator $T$ and a decomposition of the identity, and we take $D$ to be the identity operator, we recover the restricted invertibility theorem of Spielman-Srivastava.

In Section 2, we discuss some applications of our result: an improved version of Vershynin's normalized invertibility theorem and an alternative proof of the best known estimates for the proportional Dvoretzky-Rogers factorization and the Banach-Mazur distance to the cube. In Section 3, we prove our main result. Finally, in Section 4 we give a new proof of a theorem due to Kashin-Tzafriri which deals with the norm of coordinate projections of a matrix.

## 2. Applications

2.1. Normalized restricted invertibility. Our first application is an improved version of a normalized restricted invertibility principle for an arbitrary operator on $l_{2}^{n}$ and any decomposition of the identity on $\mathbb{R}^{n}$. The formulation below is essentially due to Vershynin, where the constant depending on $\varepsilon$ satisfied $c(\varepsilon) \sim \varepsilon^{-\log \varepsilon}$. We show that one can have a similar result with $c(\varepsilon)=\varepsilon$. In [16], Vershynin established a non trivial upper bound as well; he applied this to get a Dvoretzky-Rogers type lemma and to deduce information about embeddings of $l_{\infty}^{k}$ into finite dimensional spaces.

Proposition 2.1. Let $I d=\sum_{j \leqslant m} x_{j} x_{j}^{t}$ be a decomposition of the identity on $\mathbb{R}^{n}$ and let $T$ be a linear operator on $l_{2}^{n}$. For any $\varepsilon>0$ there exists $\sigma \subset\{1, \ldots, m\}$ with

$$
|\sigma| \geqslant(1-\varepsilon)^{2} \frac{\|T\|_{\mathrm{HS}}^{2}}{\|T\|_{2}^{2}}
$$

such that, for any choice of scalars $\left(a_{j}\right)_{j \in \sigma}$,

$$
\left\|\sum_{j \in \sigma} a_{j} \frac{T x_{j}}{\left\|T x_{j}\right\|_{2}}\right\|_{2} \geqslant \varepsilon\left(\sum_{j \in \sigma} a_{j}^{2}\right)^{1 / 2} .
$$

Proof. Let $U$ be the $n \times m$ matrix whose columns are the vectors $\left(T x_{j}\right)_{j \leqslant m}$. As mentioned before, we have $\|U\|_{\text {HS }}=\|T\|_{\text {HS }}$ and $\|U\|_{2}=\|T\|_{2}$. Let $D$ be the matrix with diagonal entries $\left\|T x_{j}\right\|_{2}$. Then,

$$
\|D\|_{\mathrm{HS}}=\|T\|_{\mathrm{HS}} .
$$

Now, we apply Proposition 2.1 with these $U$ and $D$ to find $\sigma$ of the desired cardinality such that

$$
s_{\min }\left(U_{\sigma} D_{\sigma}^{-1}\right) \geqslant \varepsilon .
$$

Observing that, for any $a=\left(a_{j}\right) j \in \sigma$,

$$
U_{\sigma} D_{\sigma}^{-1} a=\sum_{j \in \sigma} a_{j} \frac{T x_{j}}{\left\|T x_{j}\right\|_{2}}
$$

we conclude the proof.
2.2. Dvoretzky-Rogers factorization. We say that an $n$-dimensional normed space $X$ is in John's position if the Euclidean unit ball $B_{2}^{n}$ is the ellipsoid of maximal volume inscribed in the unit ball $B_{X}$ of $X$. In [9] (see also [1]) John characterized this position as follows: If $X$ is in John's position, then we have a decomposition $I d=\sum_{j} x_{j} x_{j}^{t}$ of the identity, where $\frac{x_{j}}{\left\|x_{j}\right\|_{2}}$ are contact points of $B_{X}$ and $B_{2}^{n}$.

The proportional Dvoretzky-Rogers factorization is given by the following statement.
Proposition 2.2. Let $X$ be a n-dimensional normed space. For any $\varepsilon>0$, there exists $k \geqslant$ $(1-\varepsilon)^{2} n$ such that the identity $i_{2, \infty}: l_{2}^{k} \longrightarrow l_{\infty}^{k}$ can be written as $i_{2, \infty}=\alpha \circ \beta$, where $\beta: l_{2}^{k} \longrightarrow X, \alpha: X \longrightarrow l_{\infty}^{k}$ and $\|\alpha\| \cdot\|\beta\| \leqslant \frac{1}{\varepsilon}$.

Proof. By changing the Euclidean structure in $\mathbb{R}^{n}$, one may assume that $X$ is in John's position. Hence, $\|\cdot\|_{X^{*}} \geqslant\|\cdot\|_{2}$ and by John's theorem one can write $I d=\sum_{j \leqslant m} x_{j} x_{j}^{t}$ where $x_{j}$ are multiples of contact points of $X$, i.e. $\left\|x_{j}\right\|_{X}=\left\|x_{j}\right\|_{X^{*}}=\left\|x_{j}\right\|_{2}$.

Now, we apply Proposition 2.1 with $T=I d$ to find $\sigma \subset\{1, \ldots, m\}$ with $|\sigma|=k \geqslant(1-\varepsilon)^{2} n$ such that:

$$
\begin{equation*}
\| \sum_{j \in \sigma} a_{j} \frac{x_{j}}{\left\|x_{j}\right\|_{2} \|_{2}} \geqslant \varepsilon\left(\sum_{j \in \sigma} a_{j}^{2}\right)^{1 / 2} \tag{1}
\end{equation*}
$$

Obviously, the vectors $\left\{x_{j}\right\}_{j \in \sigma}$ are linearly independent. Next, we define $\alpha^{\prime}: l_{1}^{k} \longrightarrow X^{*}$ by setting $\alpha^{\prime}\left(e_{j}\right)=\frac{x_{j}}{\left\|x_{j}\right\|_{2}}$, and $\beta^{\prime}: X^{*} \longrightarrow l_{2}^{k}$ by setting $\beta^{\prime}\left(\frac{x_{j}}{\left\|x_{j}\right\|_{2}}\right)=e_{j}$. Clearly, $i_{1,2}=\beta^{\prime} \circ \alpha^{\prime}$.

- Since $\left\|x_{j}\right\|_{2}=\left\|x_{j}\right\|_{X^{*}}$, we have $\left\|\alpha^{\prime}\right\|_{l_{1}^{k} \rightarrow X^{*}} \leqslant 1$.
- By (11), we have $\left\|\beta^{\prime}\right\|_{X^{*} \rightarrow l_{2}^{k}} \leqslant \frac{1}{\varepsilon}$.

Then we get the result by duality.
Remark 2.3. This form of factorization was first proven by Bourgain-Szarek (see [3]) with a weaker dependence on $\varepsilon$. After that, came the work of Szarek-Talagrand (see [14]) and of Giannopoulos (see [6]) that improved the dependence on $\varepsilon$ to $\varepsilon^{-2}$ and $\varepsilon^{-\frac{3}{2}}$ respectively. In these two works, the factorization for $i_{1,1}$ was established by combining an isomorphic variant of Sauer-Shelah's lemma with a factorization theorem of Grothendieck. Finally, Giannopoulos proved in [7] that the statement holds true with $c(\varepsilon)=\frac{c}{\varepsilon}$, which is the best known dependence on $\varepsilon$. Our Proposition recovers this result (with $c=1$ ) using a completely different method.
2.3. Banach-Mazur distance to the cube. Let $\mathbb{B M}_{n}$ denote the space of all $n$-dimensional normed spaces $X$, known as the Banach-Mazur compactum. If $X, Y$ are in $\mathbb{B M}_{n}$, we define the Banach-Mazur distance between $X$ and $Y$ as follows:

$$
\begin{aligned}
d(X, Y) & =\inf \left\{\|T\| \cdot\left\|T^{-1}\right\| \mid T \text { is an isomorphism between } \mathrm{X} \text { and } \mathrm{Y}\right\} \\
& =\inf \left\{\alpha / \beta \mid \beta B_{Y} \subset T\left(B_{X}\right) \subset \alpha B_{Y}\right\} \\
& =\inf \left\{\alpha / \beta \mid \beta\|x\|_{X} \leqslant\|T(x)\|_{Y} \leqslant \alpha\|x\|_{X}\right\} .
\end{aligned}
$$

More precisely, $\log (d)$ defines a distance (in the classical sense) in the quotient space of $\mathbb{B M}_{n}$, which is derived after we identify isometric spaces. By duality we have $d(X, Y)=d\left(X^{*}, Y^{*}\right)$.

Obtaining sharp estimates for Banach-Mazur distances is a classical topic in the asymptotic theory of finite dimensional normed spaces. The first result in this direction was a theorem of John (see [1],[9]) which provided a sharp upper bound for the distance between any $n$ dimensional normed space and the Euclidean space $\ell_{2}^{n}$. If we set $R_{2}^{n}=\max \left\{d\left(X, l_{2}^{n}\right) \mid X \in\right.$ $\left.\mathbb{B M}_{n}\right\}$, John proved that $R_{2}^{n}=\sqrt{n}$. A theorem of Gluskin (see [8]) established the existence of spaces of dimension $n$ which are at distance proportional to $n$ from each other; it follows that $l_{2}^{n}$ is at the center of $\mathbb{B M}_{n}$ and that $\mathbb{B M}_{n}$ is contained in a ball of radius $\sqrt{n}$ with respect to the Banach-Mazur distance centered at $\ell_{2}^{n}$. One natural question is to ask for an analogous estimate when $l_{2}^{n}$ is replaced by some other $l_{p}^{n}$-space, and in particular by $l_{\infty}^{n}$. This problem was first studied by Bourgain-Szarek (see [3]), then Szarek-Talagrand (see [14]) and finally Giannopou$\operatorname{los}$ (see [6]) who improved the work of Szarek-Talagrand and proved that $R_{\infty}^{n} \leqslant c n^{\frac{5}{6}}$. This is the best known result and it is probably not optimal. The Dvoretzky-Rogers factorization theorem, stated above, plays a crucial role in estimating this distance, since Proposition 2.3 allows us to find a "large subspace" of our space whose distance from $l_{2}$ and $l_{1}$ can be bounded well. For the remaining part, which will be a "small subspace", we use a trivial estimate and combine it with the previous one.

Proposition 2.4. $R_{1}^{n}=R_{\infty}^{n} \leqslant(2 n)^{\frac{5}{6}}$.

Proof. Let $X \in \mathbb{B M}_{n}$ and assume that $X^{*}$ is in John's position (by changing the Euclidean structure in $\mathbb{R}^{n}$ if necessary), i.e. $B_{2}^{n}$ is the ellipsoid of minimal volume containing $B_{X}$.
By John's theorem, one can find $x_{1}, \ldots, x_{m}$ which are multiples of contact points $\left(\left\|x_{j}\right\|_{X}=\right.$ $\left\|x_{j}\right\|_{X^{*}}=\left\|x_{j}\right\|_{2}$ ) such that $I d=\sum_{j \leqslant m} x_{j} x_{j}^{t}$. Note that $B_{2}^{n} \subset B_{X^{*}} \subset \sqrt{n} B_{2}^{n}$, and hence $\|x\|_{2} \leqslant\|x\|_{X} \leqslant \sqrt{n}\|x\|_{2}$.

Applying Proposition 2.1, one can find $\sigma \subset\{1, \ldots, m\}$ with $|\sigma|=k \geqslant(1-2 \varepsilon) n$ such that, for any choice of scalars $\left(a_{j}\right)$,

$$
\left\|\sum_{j \in \sigma} a_{j} \frac{x_{j}}{\left\|x_{j}\right\|_{2}}\right\|_{2} \geqslant \varepsilon\left(\sum_{j \in \sigma} a_{j}^{2}\right)^{1 / 2}
$$

In order to simplify notation, we assume that $\sigma=\{1, \ldots, k\}$. Let $y_{k+1}, \ldots, y_{n}$ be an orthogonal basis of $\operatorname{span}\left(x_{1}, \ldots, x_{k}\right)^{\perp}$ normalized so that $\left\|y_{j}\right\|_{X} \leqslant 1$. To do this, we require that $\left\|y_{j}\right\|_{2}=$ $\frac{1}{\sqrt{n}}$.

Define $T: l_{1}^{n} \longrightarrow X$ by setting $T\left(e_{j}\right)=\frac{x_{j}}{\left\|x_{j}\right\|_{2}}$ if $j \leqslant k$ and $T\left(e_{j}\right)=y_{j}$ for $j>k$. Let $a=\left(a_{j}\right)_{j \leqslant n} \in \mathbb{R}^{n}$ and

$$
T a=\sum_{j=1}^{k} a_{j} \frac{x_{j}}{\left\|x_{j}\right\|_{2}}+\sum_{j=k+1}^{n} a_{j} y_{j} .
$$

Then,

$$
\|a\|_{1}=\sum_{j \leqslant k}\left|a_{j}\right|+\sum_{j>k}\left|a_{j}\right| \geqslant\left\|\sum_{j \leqslant k} a_{j} \frac{x_{j}}{\left\|x_{j}\right\|_{2}}+\sum_{j>k} a_{j} y_{j}\right\|_{X} \geqslant\left\|\sum_{j \leqslant k} a_{j} \frac{x_{j}}{\left\|x_{j}\right\|_{2}}+\sum_{j>k} a_{j} y_{j}\right\|_{2} .
$$

We write

$$
\begin{aligned}
\|T a\|_{2} & \geqslant\left[\left\|\sum_{j \leqslant k} a_{j} \frac{x_{j}}{\left\|x_{j}\right\|_{2}}\right\|_{2}^{2}+\left\|\sum_{j>k} a_{j} y_{j}\right\|_{2}^{2}\right]^{\frac{1}{2}} \quad \text { by orthogonality } \\
& \geqslant\left[\varepsilon^{2} \sum_{j \leqslant k} a_{j}^{2}+\sum_{j>k} a_{j}^{2}\left\|y_{j}\right\|_{2}^{2}\right]^{\frac{1}{2}} \\
& \geqslant\left[\frac{\varepsilon^{2}}{n}\left(\sum_{j \leqslant k}\left|a_{j}\right|\right)^{2}+\frac{1}{n(n-k)}\left(\sum_{j>k}\left|a_{j}\right|\right)^{2}\right]^{\frac{1}{2}} \text { by Cauchy-Shwarz } \\
& \geqslant\left[\frac{\varepsilon^{2}}{n}\left(\sum_{j \leqslant k}\left|a_{j}\right|\right)^{2}+\frac{1}{2 \varepsilon n^{2}}\left(\sum_{j>k}\left|a_{j}\right|\right)^{2}\right]^{\frac{1}{2}} \\
& \geqslant \frac{1}{\sqrt{2}}\left[\frac{\varepsilon}{\sqrt{n}} \sum_{j \leqslant k}\left|a_{j}\right|+\frac{1}{n \sqrt{2 \varepsilon}} \sum_{j>k}\left|a_{j}\right|\right] \\
& \geqslant \frac{1}{(2 n)^{\frac{5}{6}}} \sum_{j=1}^{n}\left|a_{j}\right| \quad \text { taking } \varepsilon=(2 n)^{-\frac{1}{3}} .
\end{aligned}
$$

It follows that $(2 n)^{-\frac{5}{6}}\|a\|_{1} \leqslant\|T a\|_{X} \leqslant\|a\|_{1}$ and therefore $d\left(X, l_{1}^{n}\right) \leqslant(2 n)^{\frac{5}{6}}$ for all $X \in$ $\mathbb{B M}_{n}$.

Remark 2.5. Here we are interested in high dimensional results; this is why the constant is not that important for us. If we want an estimate for "small" dimensions, then the value of the constant becomes important. Note that a trivial estimate would be $R_{\infty}^{n} \leqslant n$. In [6], Giannopoulos
proved that $R_{\infty}^{n} \leqslant c n^{\frac{5}{6}}$ with $c=\frac{2^{\frac{7}{6}}}{(\sqrt{2}-1)^{\frac{1}{3}}} \sim 3,0116$, and thus his result becomes nontrivial when the dimension is larger than 746 . On the other hand, our result becomes nontrivial whenever the dimension is bigger than 32. Moreover, if we are interested in small dimensions, we can obtain a better result by choosing $\varepsilon$ in the last inequality in a different way: in fact we have chosen $(2 n)^{-\frac{1}{3}}$ in the asymptotic regime, otherwise one just needs to optimize $\varepsilon$ so that it satisfies $\frac{\varepsilon}{\sqrt{(1-\varepsilon)^{2} n}}=\frac{1}{n \sqrt{1-(1-\varepsilon)^{2}}}$; then our result becomes nontrivial when the dimension is larger than 15. In [15], Taschuk has also obtained an estimate for the Banach-Mazur distance to the cube of "small"-dimensional spaces. One can check that our result improves on that whenever the dimension is larger than 15 .

## 3. Proof of Theorem 1.2

Since the rank and the eigenvalues of $\left(U_{\sigma} D_{\sigma}^{-1}\right)^{t} \cdot\left(U_{\sigma} D_{\sigma}^{-1}\right)$ and $\left(U_{\sigma} D_{\sigma}^{-1}\right) \cdot\left(U_{\sigma} D_{\sigma}^{-1}\right)^{t}$ are the same, it suffices to prove that $\left(U_{\sigma} D_{\sigma}^{-1}\right) \cdot\left(U_{\sigma} D_{\sigma}^{-1}\right)^{t}$ has rank equal to $k=|\sigma|$ and its smallest positive eigenvalue is greater than $\varepsilon^{2} \frac{\|U\|_{\mathrm{HS}}^{2}}{\|D\|_{\mathrm{HS}}}$. Note that

$$
\left(U_{\sigma} D_{\sigma}^{-1}\right) \cdot\left(U_{\sigma} D_{\sigma}^{-1}\right)^{t}=\sum_{j \in \sigma}\left(U D_{\sigma}^{-1} e_{j}\right) \cdot\left(U D_{\sigma}^{-1} e_{j}\right)^{t}=\sum_{j \in \sigma}\left(\frac{U e_{j}}{\alpha_{j}}\right) \cdot\left(\frac{U e_{j}}{\alpha_{j}}\right)^{t}
$$

We are going to construct the matrix $A_{k}=\sum_{j \in \sigma}\left(U D_{\sigma}^{-1} e_{j}\right) \cdot\left(U D_{\sigma}^{-1} e_{j}\right)^{t}$ by iteration. We begin by setting $A_{0}=0$ and at each step we will be adding a rank one matrix $\left(\left(\frac{U e_{j}}{\alpha_{j}}\right) \cdot\left(\frac{U e_{j}}{\alpha_{j}}\right)^{t}\right.$ for a suitable $j$ ) which will give a new positive eigenvalue. This will guarantee that the vector $U D_{\sigma}^{-1} e_{j}$ chosen in each step is linearly independent from the previous ones.

If $A$ and $B$ are symmetric matrices, we write $A \preceq B$ if $B-A$ is a positive semidefinite matrix. Recall the Sherman-Morrison Formula which will be needed in the proof. For any invertible matrix $A$ and any vector $v$ we have

$$
\left(A+v \cdot v^{t}\right)^{-1}=A^{-1}-\frac{A^{-1} v \cdot v^{t} A^{-1}}{1+v^{t} A^{-1} v}
$$

We will also apply the following lemma which appears as Lemma 6.3 in [13]:
Lemma 3.1. Suppose that $A \succeq 0$ has $q$ nonzero eigenvalues, all greater than $b^{\prime}>0$. If $v \neq 0$ and

$$
\begin{equation*}
v^{t}\left(A-b^{\prime} I\right)^{-1} v<-1 \tag{2}
\end{equation*}
$$

then $A+v v^{t}$ has $q+1$ nonzero eigenvalues, all greater than $b^{\prime}$.
The proof of the lemma is simple and makes use of the Sherman-Morrison formula.
For any symmetric matrix $A$ and any $b>0$, we define

$$
\phi_{b}(A)=\operatorname{Tr}\left(U^{t}(A-b I)^{-1} U\right)
$$

as the potential corresponding to the barrier $b$.
At each step $l$, the matrix already constructed is denoted by $A_{l}$ and the barrier by $b_{l}$ and to simplify the notations, we will use $\phi_{l}$ for the potential $\phi_{b_{l}}$. $A_{l}$ has $l$ nonzero eigenvalues all greater than $b_{l}$. As mentioned before, we will try to construct $A_{l+1}$ by adding a rank one matrix $v \cdot v^{t}$ to $A_{l}$ so that $A_{l+1}$ has $l+1$ nonzero eigenvalues all greater than $b_{l+1}=b_{l}-\delta$ and
$\phi_{l+1}\left(A_{l+1}\right) \leqslant \phi_{l}\left(A_{l}\right)$. Note that

$$
\begin{aligned}
\phi_{l+1}\left(A_{l+1}\right) & =\operatorname{Tr}\left(U^{t}\left(A_{l}+v v^{t}-b_{l+1} I\right)^{-1} U\right) \\
& =\operatorname{Tr}\left(U^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U\right)-\operatorname{Tr}\left(\frac{U^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U}{1+v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v}\right) \\
& =\phi_{l+1}\left(A_{l}\right)-\frac{v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U U^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v}{1+v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v} .
\end{aligned}
$$

So, in order to have $\phi_{l+1}\left(A_{l+1}\right) \leqslant \phi_{l}\left(A_{l}\right)$, we must choose a vector $v$ verifying

$$
\begin{equation*}
-\frac{v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U U^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v}{1+v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v} \leqslant \phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right) . \tag{3}
\end{equation*}
$$

Since $v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U U^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v$ and $\left.\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)\right)$ are positive, choosing $v$ verifying conditions (2) and (3) is equivalent to choosing $v$ which satisfies the following:
(4) $v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U U^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v \leqslant\left(\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)\right)\left(-1-v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v\right)$.

Since $U U^{t} \preceq\|U\|_{2}^{2} I d$ and $\left(A_{l}-b_{l+1} I\right)^{-1}$ is symmetric, it is sufficient to choose $v$ so that

$$
\begin{equation*}
v^{t}\left(A_{l}-b_{l+1} I\right)^{-2} v \leqslant \frac{1}{\|U\|_{2}^{2}}\left(\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)\right)\left(-1-v^{t}\left(A_{l}-b_{l+1} I\right)^{-1} v\right) \tag{5}
\end{equation*}
$$

Recall the notation $I_{D}:=\left\{j \leqslant m \mid \alpha_{j} \neq 0\right\}$ where $\left(\alpha_{j}\right)_{j \leqslant m}$ are the diagonal entries of $D$. Since we have assumed that $\operatorname{Ker}(D) \subset \operatorname{Ker}(U)$, we have

$$
\|U\|_{\mathrm{HS}}^{2}=\sum_{j \leqslant m}\left\|U e_{j}\right\|_{2}^{2}=\sum_{j \in I_{D}}\left\|U e_{j}\right\|_{2}^{2} \leqslant\left|I_{D}\right| \cdot\|U\|_{2}^{2}
$$

and thus $\left|I_{D}\right| \geqslant \frac{\|U\|_{\mathrm{HS}}^{2}}{\|U\|_{2}^{2}}$. At each step, we will select a vector $v$ satisfying (5) among $\left(\frac{U e_{j}}{\alpha_{j}}\right)_{j \in I_{D}}$. Our task therefore is to find $j \in I_{D}$ such that

$$
\begin{equation*}
\left(U e_{j}\right)^{t}\left(A_{l}-b_{l+1} I\right)^{-2} U e_{j} \leqslant \frac{\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)}{\|U\|_{2}^{2}}\left(-\alpha_{j}^{2}-\left(U e_{j}\right)^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U e_{j}\right) \tag{6}
\end{equation*}
$$

The existence of such a $j \in I_{D}$ is guaranteed by the fact that condition (6) holds true if we take the sum over all $\left(\frac{U e_{j}}{\alpha_{j}}\right)_{j \in D}$. The hypothesis $\operatorname{Ker}(D) \subset \operatorname{Ker}(U)$ implies that:

- $\sum_{j \in I_{D}}\left(U e_{j}\right)^{t}\left(A_{l}-b_{l+1} I\right)^{-2} U e_{j}=\operatorname{Tr}\left(U^{t}\left(A_{l}-b_{l+1} I\right)^{-2} U\right)$,
- $\sum_{j \in I_{D}}\left(U e_{j}\right)^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U e_{j}=\operatorname{Tr}\left(U^{t}\left(A_{l}-b_{l+1} I\right)^{-1} U\right)$.

Therefore it is enough to prove that, at each step, one has

$$
\begin{equation*}
\operatorname{Tr}\left(U^{t}\left(A_{l}-b_{l+1} I\right)^{-2} U\right) \leqslant \frac{\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)}{\|U\|_{2}^{2}}\left(-\|D\|_{\mathrm{HS}}^{2}-\phi_{l+1}\left(A_{l}\right)\right) . \tag{7}
\end{equation*}
$$

The rest of the proof is similar to the one in [13]. One just needs to replace $m$ by $\|D\|_{\mathrm{HS}}^{2}$. For the sake of completeness, we include the proof. The next lemma will determine the conditions required at each step in order to prove (7).
Lemma 3.2. Suppose that $A_{l}$ has $l$ nonzero eigenvalues all greater than $b_{l}$, and write $Z$ for the orthogonal projection onto the kernel of $A_{l}$. If

$$
\begin{equation*}
\phi_{l}\left(A_{l}\right) \leqslant-\|D\|_{\mathrm{HS}}^{2}-\frac{\|U\|_{2}^{2}}{\delta} \tag{8}
\end{equation*}
$$

and

$$
\begin{equation*}
0<\delta<b_{l} \leqslant \delta \frac{\|Z U\|_{\mathrm{HS}}^{2}}{\|U\|_{2}^{2}} \tag{9}
\end{equation*}
$$

then there exists $i \in I_{D}$ such that $A_{l+1}:=A_{l}+\left(\frac{U e_{i}}{\alpha_{i}}\right) \cdot\left(\frac{U e_{i}}{\alpha_{i}}\right)^{t}$ has $l+1$ nonzero eigenvalues all greater than $b_{l+1}:=b_{l}-\delta$ and $\phi_{l+1}\left(A_{l+1}\right) \leqslant \phi_{l}\left(A_{l}\right)$.
Proof. As mentioned before, it is enough to prove inequality (7). We set $\Delta_{l}:=\phi_{l}\left(A_{l}\right)-$ $\phi_{l+1}\left(A_{l+1}\right)$. By (8), we get

$$
\phi_{l+1}\left(A_{l}\right) \leqslant-\|D\|_{\mathrm{HS}}^{2}-\frac{\|U\|_{2}^{2}}{\delta}-\Delta_{l} .
$$

Inserting this in (7), we see that it is sufficient to prove the following inequality:

$$
\begin{equation*}
\operatorname{Tr}\left(U^{t}\left(A_{l}-b_{l+1} I\right)^{-2} U\right) \leqslant \Delta_{l}\left(\frac{\Delta_{l}}{\|U\|_{2}^{2}}+\frac{1}{\delta}\right) \tag{10}
\end{equation*}
$$

Now, denote by $P$ the orthogonal projection onto the image of $A_{l}$. We set

$$
\phi_{l}^{P}\left(A_{l}\right):=\operatorname{Tr}\left(U^{t} P\left(A_{l}-b_{l} I\right)^{-1} P U\right) \quad \text { and } \quad \Delta_{l}^{P}:=\phi_{l}^{P}\left(A_{l}\right)-\phi_{l+1}^{P}\left(A_{l}\right)
$$

and use similar notation for $Z$. Since $P, Z$ and $A_{l}$ commute, one can write

$$
\Delta_{l}=\Delta_{l}^{P}+\Delta_{l}^{Z} \quad \text { and } \quad \phi_{l}\left(A_{l}\right)=\phi_{l}^{P}\left(A_{l}\right)+\phi_{l}^{Z}\left(A_{l}\right)
$$

Note that:

$$
\begin{aligned}
\left(A_{l}-b_{l} I\right)^{-1}-\left(A_{l}-b_{l+1} I\right)^{-1} & =\left(A_{l}-b_{l} I\right)^{-1}\left(b_{l} I-A_{l}+A_{l}-b_{l+1} I\right)\left(A_{l}-b_{l+1} I\right)^{-1} \\
& =\delta\left(A_{l}-b_{l} I\right)^{-1}\left(A_{l}-b_{l+1} I\right)^{-1}
\end{aligned}
$$

and since $P\left(A_{l}-b_{l} I\right)^{-1} P$ and $P\left(A_{l}-b_{l+1} I\right)^{-1} P$ are positive semidefinite, we have:

$$
U^{t} P\left(A_{l}-b_{l} I\right)^{-1} P U-U^{t} P\left(A_{l}-b_{l+1} I\right)^{-1} P U \succeq \delta U^{t} P\left(A_{l}-b_{l+1} I\right)^{-2} P U
$$

Inserting this in (10), it is enough to prove that:

$$
\operatorname{Tr}\left(U^{t} Z\left(A_{l}-b_{l+1} I\right)^{-2} Z U\right) \leqslant \Delta_{l}\left(\frac{\Delta_{l}}{\|U\|_{2}^{2}}+\frac{1}{\delta}\right)-\frac{\Delta_{l}^{P}}{\delta}
$$

Since $A_{l} Z=0$, we have:

- $\operatorname{Tr}\left(U^{t} Z\left(A_{l}-b_{l+1} I\right)^{-2} Z U\right)=\frac{\|Z U\|_{\text {HS }}^{2}}{b_{l+1}^{2}}$ and
- $\Delta_{l}^{Z}=-\frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l}}+\frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l+1}}=\delta \frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l} b_{l+1}}$,
so taking into account the fact that $\Delta_{l} \geqslant \Delta_{l}^{Z} \geqslant 0$, it remains to prove the following:

$$
\begin{equation*}
\frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l+1}^{2}} \leqslant \delta^{2} \frac{\|Z U\|_{\mathrm{HS}}^{4}}{\|U\|_{2}^{2} b_{l}^{2} b_{l+1}^{2}}+\frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l} b_{l+1}} . \tag{11}
\end{equation*}
$$

By Hypothesis (9), this last inequality follows by

$$
\begin{equation*}
\frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l+1}^{2}} \leqslant \delta \frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l} b_{l+1}^{2}}+\frac{\|Z U\|_{\mathrm{HS}}^{2}}{b_{l} b_{l+1}} \tag{12}
\end{equation*}
$$

which is trivially true since $b_{l+1}=b_{l}-\delta$.
We are now able to complete the proof of Theorem 1.2, To this end, we must verify that conditions (8) and (9) hold at each step. At the beginning we have $A_{0}=0$ and $Z=I d$, so we must choose a barrier $b_{0}$ such that:

$$
\begin{equation*}
-\frac{\|U\|_{\mathrm{HS}}^{2}}{b_{0}} \leqslant-\|D\|_{\mathrm{HS}}^{2}-\frac{\|U\|_{2}^{2}}{\delta} \tag{13}
\end{equation*}
$$

and

$$
\begin{equation*}
b_{0} \leqslant \delta \frac{\|U\|_{\mathrm{HS}}^{2}}{\|U\|_{2}^{2}} \tag{14}
\end{equation*}
$$

We choose

$$
b_{0}:=\varepsilon \frac{\|U\|_{\mathrm{HS}}^{2}}{\|D\|_{\mathrm{HS}}^{2}} \quad \text { and } \quad \delta:=\frac{\varepsilon}{1-\varepsilon} \frac{\|U\|_{2}^{2}}{\|D\|_{\mathrm{HS}}^{2}}
$$

and we note that (13) and (14) are verified. Also, at each step (8) holds because $\phi_{l+1}\left(A_{l+1}\right) \leqslant$ $\phi_{l}\left(A_{l}\right)$. Since $\|Z U\|_{\text {HS }}^{2}$ decreases at each step by at most $\|U\|_{2}^{2}$, the right-hand side of (9) decreases by at most $\delta$, and therefore (9) holds once we replace $b_{l}$ by $b_{l}-\delta$.

Finally note that, after $k=(1-\varepsilon)^{2} \frac{\|U\|_{\mathrm{HS}}^{2}}{\|U\|_{2}^{2}}$ steps, the barrier will be

$$
b_{k}=b_{0}-k \delta=\varepsilon^{2} \frac{\|U\|_{\mathrm{HS}}^{2}}{\|D\|_{\mathrm{HS}}^{2}}
$$

This completes the proof.

## 4. PRojection on coordinate subspaces

Given an $n \times m$ rectangular matrix, a theorem of Kashin-Tzafriri (see [10]) allows us to find a coordinate projection of arbitrary size (say $\lambda m$ with $\lambda<1$ ) and to estimate its norm in terms of the norm of the matrix and $\lambda$. The precise formulation of the result (see [16] for another formulation and proof) is the following.

Theorem 4.1 (Kashin-Tzafriri). Let $U$ be an $n \times m$ matrix. Fix $\lambda$ with $1 / m \leqslant \lambda \leqslant \frac{1}{4}$. Then, there exists a subset $\nu$ of $\{1, \ldots, m\}$ of cardinality $|\nu| \geqslant \lambda m$ such that

$$
\left\|U_{\nu}\right\| \leqslant c\left(\sqrt{\lambda}\|U\|_{2}+\frac{\|U\|_{\mathrm{HS}}}{\sqrt{m}}\right)
$$

where $U_{\nu}=U P_{\nu}$ and $P_{\nu}$ denotes the coordinate projection onto $\mathbb{R}^{\nu}$.
The proof of this theorem uses a standard selection argument combined with a factorization theorem of Grothendieck. We will give an alternative proof of this result which exploits the method introduced by Batson-Spielman-Srivastava in [2]. This allows us to slightly improve Theorem 4.1 by obtaining an estimate for all possible values of the rank of the projection with the constant $c$ replaced by $\frac{\sqrt{2}}{\sqrt{1-\lambda}}$.

Theorem 4.2. Let $U$ be an $n \times m$ matrix and let $0<\lambda \leqslant \eta<1$. Then, there exists $\sigma \subset$ $\{1, \ldots, m\}$ with $|\sigma|=k \geqslant \lambda m$ such that

$$
\left\|U_{\sigma}\right\|_{2} \leqslant \frac{1}{\sqrt{1-\lambda}}\left(\sqrt{\lambda+\eta}\|U\|_{2}+\sqrt{1+\frac{\lambda}{\eta}} \frac{\|U\|_{\mathrm{HS}}}{\sqrt{m}}\right)
$$

In particular,

$$
\left\|U_{\sigma}\right\|_{2} \leqslant \frac{\sqrt{2}}{\sqrt{1-\lambda}}\left(\sqrt{\lambda}\|U\|_{2}+\frac{\|U\|_{\mathrm{HS}}}{\sqrt{m}}\right)
$$

where $U_{\sigma}$ denotes the selection of the columns of $U$ with indices in $\sigma$.
Proof. We denote by $\left(e_{j}\right)_{j \leqslant m}$ the canonical basis of $\mathbb{R}^{m}$. Since

$$
U_{\sigma} \cdot U_{\sigma}^{t}=\sum_{j \leqslant \sigma}\left(U e_{j}\right) \cdot\left(U e_{j}\right)^{t}
$$

our problem reduces to the question of estimating the largest eigenvalue of this sum of rank one matrices. We will follow the same procedure as in the proof of the restricted invertibility theorem: at each step, we would like to add a column of the original matrix and then study the evolution of the largest eigenvalue. However, it will be convenient for us to add suitable multiples of the columns of $U$ in order to construct the $k$-th matrix; for each $k$ we will choose a subset $\sigma_{k}$ of cardinality $\left|\sigma_{k}\right|=k$ and consider the matrix $A_{k}=\sum_{j \in \sigma_{k}} s_{j}\left(U e_{j}\right) \cdot\left(U e_{j}\right)^{t}$ where
$\left(s_{j}\right)_{j \in \sigma}$ will be positive numbers which will be suitably chosen. At the step $l$, the barrier will be denoted by $b_{l}$, namely the eigenvalues of $A_{l}$ will be all smaller than $b_{l}$. The corresponding potential is $\phi_{l}\left(A_{l}\right):=\operatorname{Tr}\left(U^{t}\left(b_{l} I-A_{l}\right)^{-1} U\right)$. We set $A_{0}=0$, while $b_{0}$ will be determined later.

As we did before, at each step the value of the potential $\phi_{l}\left(A_{l}\right)$ will decrease so that we can continue the iteration, while the value of the barrier will increase by a constant $\delta$, i.e. $b_{l+1}=b_{l}+\delta$. We will use a lemma which appears as Lemma 3.4 in [13]. We state it here in the notation introduced above.

Lemma 4.3. Assume that $\lambda_{\max }\left(A_{l}\right) \leqslant b_{l}$. Let $v$ be a vector in $\mathbb{R}^{n}$ satisfying

$$
\mathbb{V}_{l}(v):=\frac{v^{t}\left(b_{l+1} I-A_{l}\right)^{-2} v}{\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)}\|U\|_{2}+v^{t}\left(b_{l+1} I-A_{l}\right)^{-1} v \leqslant \frac{1}{s} .
$$

Then, if we define $A_{l+1}=A_{l}+\operatorname{svv}^{t}$ we have

$$
\lambda_{\max }\left(A_{l+1}\right) \leqslant b_{l+1} \quad \text { and } \quad \phi_{l+1}\left(A_{l+1}\right) \leqslant \phi_{l}\left(A_{l}\right)
$$

We write $\alpha$ for the initial potential, i.e. $\alpha=\frac{\|U\|_{\mathrm{HS}}^{2}}{b_{0}}$. Suppose that $A_{l}=\sum_{j \in \sigma_{l}} s_{j}\left(U e_{j}\right) \cdot\left(U e_{j}\right)^{t}$ is constructed so that $\phi_{l}\left(A_{l}\right) \leqslant \phi_{l-1}\left(A_{l-1}\right) \leqslant \alpha$ and $\lambda_{\max }\left(A_{l}\right) \leqslant b_{l}$. We will now use Lemma4.3 in order to construct $A_{l+1}$. To this end, we must find a vector $U e_{j}$ not chosen before and a scalar $s_{l+1}$ so that $\mathbb{V}_{l}\left(U e_{j}\right) \leqslant \frac{1}{s_{l+1}}$, and then use the lemma. We choose $j$ by calculating the sum of $\mathbb{V}_{l}\left(U e_{j}\right)$ over all $j \notin \sigma_{l}$. Since $\left(b_{l} I-A_{l}\right)^{-1}$ and $\left(b_{l+1} I-A_{l}\right)^{-1}$ are positive semidefinite, one can easily check that

$$
\left(b_{l} I-A_{l}\right)^{-1}-\left(b_{l+1} I-A_{l}\right)^{-1} \succeq \delta\left(b_{l+1} I-A_{l}\right)^{-2}
$$

Therefore,

$$
\operatorname{Tr}\left(U^{t}\left(b_{l+1} I-A_{l}\right)^{-2} U\right) \leqslant \frac{1}{\delta}\left(\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)\right)
$$

It follows that

$$
\begin{aligned}
\sum_{j \notin \sigma_{l}} \mathbb{V}_{l}\left(U e_{j}\right) & \leqslant \sum_{j \leqslant m} \mathbb{V}_{l}\left(U e_{j}\right)=\frac{\operatorname{Tr}\left(U^{t}\left(b_{l+1} I-A_{l}\right)^{-2} U\right)}{\phi_{l}\left(A_{l}\right)-\phi_{l+1}\left(A_{l}\right)}\|U\|_{2}^{2}+\phi_{l+1}\left(A_{l}\right) \\
& \leqslant \frac{\|U\|_{2}^{2}}{\delta}+\alpha,
\end{aligned}
$$

and therefore one can find $i \notin \sigma_{l}$ such that

$$
\mathbb{V}_{l}\left(U e_{i}\right) \leqslant \frac{1}{\left|\sigma_{l}^{c}\right|}\left(\frac{\|U\|_{2}^{2}}{\delta}+\alpha\right) \leqslant \frac{1}{\left|\sigma_{k}^{c}\right|}\left(\frac{\|U\|_{2}^{2}}{\delta}+\alpha\right)
$$

where $k$ is the maximum number of steps (which is in our case $\lambda m$ ).
 take $A_{l+1}=A_{l}+s\left(U e_{i}\right) \cdot\left(U e_{i}\right)^{t}$. After $k=\lambda m$ steps, we get $\sigma=\sigma_{k}$ such that

$$
\begin{aligned}
\lambda_{\max }\left(\sum_{j \in \sigma_{k}}\left(U e_{j}\right) \cdot\left(U e_{j}\right)^{t}\right) & \leqslant \frac{1}{s}\left(b_{0}+k \delta\right)=\frac{\alpha+\frac{\|U\|_{2}^{2}}{\delta}}{(1-\lambda) m}\left(b_{0}+k \delta\right) \\
& =\frac{1}{1-\lambda}\left[\frac{\|U\|_{\mathrm{HS}}^{2}}{m}+\lambda\|U\|_{2}^{2}+\lambda\|U\|_{\mathrm{HS}}^{2} \frac{\delta}{b_{0}}+\frac{\|U\|_{2}^{2}}{m} \frac{b_{0}}{\delta}\right]
\end{aligned}
$$

The result follows by taking $b_{0}=\eta m \delta$.
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