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# Moments Tensors, Hilbert's Identity, and $k$-wise Uncorrelated Random Variables 

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#### Abstract

In this paper, we introduce a notion to be called $k$-wise uncorrelated random variables, which is similar but not identical to the so-called $k$-wise independent random variables in the literature. We show how to construct $k$-wise uncorrelated random variables by a simple procedure. The constructed random variables can be applied, e.g., to express the quartic polynomial $\left(x^{\mathrm{T}} Q x\right)^{2}$, where $Q$ is an $n \times n$ positive semidefinite matrix, by a sum of fourth powered polynomial terms, known as Hilbert's identity. By virtue of the proposed construction, the number of required terms is no more than $2 n^{4}+n$. This implies that it is possible to find a $\left(2 n^{4}+n\right)$-point distribution whose fourth moments tensor is exactly the symmetrization of $Q \otimes Q$. Moreover, we prove that the number of required fourth powered polynomial terms to express $\left(x^{\mathrm{T}} Q x\right)^{2}$ is at least $n(n+1) / 2$. The result is applied to prove that computing the matrix $2 \mapsto 4$ norm is NP-hard. Extensions of the results to complex random variables are discussed as well.


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1. Introduction. Given an $n$-dimensional random vector $\xi=\left(\xi_{1}, \xi_{2}, \ldots, \xi_{n}\right)^{\mathrm{T}}$ with joint density function $p(\cdot)$, let us denote the $n$-dimensional $d$-th order tensor $\mathscr{F}$ to be the $d$-th order moments tensor associated with $\xi$ as follows:

$$
\mathscr{F}_{i_{1}, i_{2}, \ldots, i_{d}}=\mathrm{E}\left[\prod_{k=1}^{d} \xi_{i_{k}}\right]=\int_{\mathbb{R}^{n}} \prod_{k=1}^{d} u_{i_{k}} p(u) d u \quad \forall 1 \leq i_{1}, i_{2}, \ldots, i_{d} \leq n
$$

or equivalently,

$$
\mathscr{F}=\int_{\mathbb{R}^{n}} \underbrace{u \otimes u \otimes \cdots \otimes u}_{d} p(u) d u .
$$

Because tensor $\mathscr{F}$ is in a finite dimensional space, by Carathéodory's theorem (Carathéodory [6]), it can be further rewritten as a sum of finite "rank one" terms, i.e., there exist $t$ vectors $b^{1}, b^{2}, \ldots, b^{t}$ such that

$$
\begin{equation*}
\mathscr{F}=\sum_{i=1}^{t} \underbrace{b^{i} \otimes b^{i} \otimes \cdots \otimes b^{i}}_{d} \tag{1}
\end{equation*}
$$

An immediate consequence of the above construction is that $\mathscr{F}$ is supersymmetric, meaning that its component is invariant under permutation of the indices. For instance, the second order moments tensor can be easily derived from its covariance matrix, which is naturally symmetric and positive semidefinite. Indeed, thanks to the formulation (1), any $2 d$-th order moments tensor is always positive semidefinite, in other words, the homogeneous polynomial function induced by this tensor is always nonnegative, i.e.,

$$
f(x)=\mathscr{F}(\underbrace{x, x, \ldots, x}_{2 d}):=\sum_{1 \leq i_{1}, i_{2}, \ldots, i_{2 d} \leq n} \mathscr{F}_{i_{1}, i_{2}, \ldots, i_{2 d}} \prod_{k=1}^{2 d} x_{i_{k}}=\sum_{i=1}^{t}\left(\left(b^{i}\right)^{\mathrm{T}} x\right)^{2 d} \geq 0 .
$$

However, the term "nonnegativity" can be ambiguous in the case of higher order tensors. In our recent paper, Jiang et al. [11], this issue was particularly addressed. We shall only note here that the $2 d$-th moments tensors form a specific nonnegative convex cone, whose membership query is a hard problem in general (see Jiang et al. [11]). It is therefore interesting to know what kind of tensors are contained in this cone. For instance, one may wonder if the supersymmetric tensor associated with the polynomial $\left(x^{\mathrm{T}} x\right)^{2}$, which is clearly nonnegative, is a fourth order moments tensor or not. Interestingly, the answer is yes, because of a result of Hilbert [10], who showed that it is possible to express $\left(x^{\mathrm{T}} x\right)^{d}$ as $\sum_{i=1}^{t}\left(x^{\mathrm{T}} a^{i}\right)^{2 d}$. As a consequence, the polynomial $\left(x^{\mathrm{T}} x\right)^{2}$ (the case $d=2$ ) can be viewed as $\mathrm{E}\left[\xi^{\mathrm{T}} x\right]^{4}$ where $\xi$ is a random vector, taking value $t^{1 / 4} a^{i}$ with probability $1 / t$. Therefore, $\operatorname{sym}(I \otimes I)$ with $I$ being the identity matrix is a fourth moments tensor, where the symmetrization mapping "sym" turns a given tensor into a supersymmetric one by making the entries with the same set of indices all the same (taking the value of the average).

Apart from the above example, there are several other representations for general $2 d$-th moments tensor other than (1). For example, with the help of Hilbert's identity (Barvinok [4]), we can easily verify that

with $A \succeq 0$ also belongs to $2 d$-th moments cone. Specifically, one can find vectors $a^{1}, a^{2}, \ldots, a^{t}$ such that

$$
\begin{equation*}
\operatorname{sym}(\underbrace{A \otimes A \otimes \cdots \otimes A}_{d})=\sum_{i=1}^{t} \underbrace{a^{i} \otimes a^{i} \otimes \cdots \otimes a^{i}}_{2 d} \tag{2}
\end{equation*}
$$

On the other hand, by letting the order of the tensor be $2 d$ and $A^{i}=b^{i} \otimes b^{i}=b^{i}\left(b^{i}\right)^{\mathrm{T}}$ in (1), we have

$$
\begin{equation*}
\mathscr{F}=\sum_{i=1}^{t} \underbrace{b^{i} \otimes b^{i} \otimes \cdots \otimes b^{i}}_{2 d}=\sum_{i=1}^{t} \operatorname{sym}(\underbrace{A^{i} \otimes A^{i} \otimes \cdots \otimes A^{i}}_{d}) \quad \text { with } A^{i} \succeq 0 \text { and } \operatorname{rank}\left(A^{i}\right)=1 \tag{3}
\end{equation*}
$$

This implies that the rank-one constraint is redundant in terms of requiring $\mathscr{F}$ to be a $2 d$-th moments tensor in (3).

In general, such decomposition of (2) is not unique. For example, one may verify that

$$
\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{2}=\frac{1}{3} \sum_{i=1}^{3} x_{i}^{4}+\frac{1}{3} \sum_{1 \leq i<j \leq 3} \sum_{\beta_{j}= \pm 1}\left(x_{i}+\beta_{j} x_{j}\right)^{4}=\frac{2}{3} \sum_{i=1}^{3} x_{i}^{4}+\frac{1}{3} \sum_{\beta_{2}, \beta_{3}= \pm 1}\left(x_{1}+\beta_{2} x_{2}+\beta_{3} x_{3}\right)^{4},
$$

which leads to two different representations of the tensor $\operatorname{sym}\left(I_{3} \otimes I_{3}\right)$. An interesting question is to find a succinct (preferably the shortest) representation among all the different representations, including the one from Hilbert's decomposition. However, from the original Hilbert's construction, the representation on the right-hand side of (2) is exponential in $n$. By Carathéodory's theorem, there exists a decomposition such that the value of $t$ in (2) is no more than $\binom{n+2 d-1}{2 d}+1$. Unfortunately, Carathéodory's theorem is nonconstructive. This motivates us to construct a polynomial-size representation, i.e., $t=O\left(n^{k}\right)$ for some constant $k$ in (2).

One contribution of this paper is to give a "short" (polynomial-size) representation for Hilbert's identity when $d=2$. In fact, we also prove the number of terms for any representation can never be less than $n(n+1) / 2$. An application of this polynomial-size representation will be discussed. Toward this end, let us first introduce the new notion of $k$-wise uncorrelated random variables, which may appear to be completely unrelated to the discussion of Hilbert's identity at first glance.

Definition 1.1 ( $k$-wise Uncorrelation). A set of random variables $\left\{\xi_{1}, \xi_{2}, \ldots, \xi_{n}\right\}$ is called $k$-wise uncorrelated if

$$
\mathrm{E}\left[\prod_{j=1}^{n} \xi_{j}^{p_{j}}\right]=\prod_{j=1}^{n} \mathrm{E}\left[\xi_{j}^{p_{j}}\right] \quad \forall p_{1}, p_{2}, \ldots, p_{n} \in \mathbb{Z}_{+} \quad \text { with } \quad \sum_{i=1}^{n} p_{i}=k
$$

For instance, if $\xi_{1}, \xi_{2}, \ldots, \xi_{n}$ are independent and identically distributed (i.i.d.) random variables with finite supporting set $|\Delta|=q$, then they are $k$-wise uncorrelated. However the size of its corresponding sample space is $q^{n}$, which is exponential in $n$. It turns out that reducing the sample space while keeping the $k$-wise uncorrelation structure can be of great importance in many applications. For example, our result shows that the polynomialsize representation (2) can be obtained by finding $k$-wise uncorrelated random variables with polynomial-size sample space. Before addressing the issue of finding such random variables, we shall first discuss a related notion known as the $k$-wise independence.

Definition 1.2 ( $k$-wise Independence). A set of random variables $\Xi=\left\{\xi_{1}, \xi_{2}, \ldots, \xi_{n}\right\}$ with each taking values on the set $\Delta=\left\{\delta_{1}, \delta_{2}, \ldots, \delta_{q}\right\}$ is called $k$-wise independent, if any $k$ different random variables $\xi_{i_{1}}, \xi_{i_{2}}, \ldots, \xi_{i_{k}}$ of $\Xi$ are independent, i.e.,

$$
\operatorname{Prob}\left\{\xi_{i_{1}}=\delta_{i_{1}}, \xi_{i_{2}}=\delta_{i_{2}}, \ldots, \xi_{i_{k}}=\delta_{i_{k}}\right\}=\prod_{j=1}^{k} \operatorname{Prob}\left\{\xi_{i_{j}}=\delta_{i_{j}}\right\} \quad \forall \delta_{i_{j}} \in \Delta, \quad j=1,2, \ldots, k
$$

Note that when $k=2$, $k$-wise independence is usually called pair-wise independence. Since the $1980 \mathrm{~s}, k$-wise independence has been a popular topic in theoretical computer science. Essentially, working with $k$-wise independence (instead of the full independence) means that one can reduce the size of the sample space in question. In many cases, this feature is crucial. For instance, when $\Delta=\{0,1\}$ and $\operatorname{Prob}\left\{\xi_{1}=0\right\}=\operatorname{Prob}\left\{\xi_{1}=1\right\}=\frac{1}{2}$, Alon et al. [2] constructed a sample space of size being approximately $n^{k / 2}$. For the same $\Delta$, when $\xi_{1}, \xi_{2}, \ldots, \xi_{n}$ are independent but not identical, Karloff and Mansour [13] proved that the size of sample space can be upper bounded by $O\left(n^{k}\right)$. In the case of $\Delta=\{0,1, \ldots, q-1\}$ with $q$ being a prime number, the total number of random variables being $k$-wise independent are quite restricted. For given $k<q$, Joffe [12] showed that there are up to $q+1$ random variables form a $k$-wise independent set and the size of the sample space is $q^{k}$.

Clearly, $k$-wise independence implies $k$-wise uncorrelation. Therefore, we may apply the existing results of $k$-wise independence to get $k$-wise uncorrelated random variables. However, the aforementioned constructions of $k$-wise independent random variables heavily depend on the structure of $\Delta$ (e.g., it requires that $|\Delta|=2$ or $k<|\Delta|$ ). Moreover, the construction of $k$-wise independent random variables is typically complicated and technically involved (see Karloff and Mansour [13]). In fact, for certain problems (e.g., polynomial-size representation of Hilbert's identity in this case), we only need the random variables to be $k$-wise uncorrelated. Therefore, in this paper we propose a tailor-made simple construction that suits the structure of $k$-wise uncorrelated random variables. As we shall see later, our approach can handle the more general support set:

$$
\begin{equation*}
\Delta_{q}:=\left\{1, \omega_{q}, \ldots, \omega_{q}^{q-1}\right\} \quad \text { with } \omega_{q}=e^{\mathbf{i}((2 \pi) / q)}=\cos \frac{2 \pi}{q}+\mathbf{i} \sin \frac{2 \pi}{q} \text { and } q \text { is prime, } \tag{4}
\end{equation*}
$$

and $k$ can be any parameter.
Conceptually, our approach is rather generic: the $k$-wise uncorrelated random variables are constructed based only on the product of a small set of i.i.d. random variables with their powers; the sample space would be polynomial size if the number of such i.i.d. random variables is $O(\log n)$. Consequently, we not only find polynomial-size representation for the fourth moments tensor in the form of $\operatorname{sym}(A \otimes A)$, but also for complex $2^{d} q$-th moments tensor. As an application, this construction can be used to prove that the matrix $2 \mapsto 4$ norm problem (Bhaskara and Vijayaraghavan [5]), whose complexity was previously unknown, ${ }^{1}$ is actually NP-hard.

The rest of this paper is organized as follows. In $\S 2$ we introduce Hilbert's identity and its connections to $2 d$ th moments tensor. Then, in $\S 3$ we present a randomized algorithm, as well as a deterministic one, to construct $k$-wise uncorrelated random variables. As a result, we find polynomial-size representation of fourth moments tensor and complex $2^{d} q$-th moments tensor in $\S 4$. In $\S 5$, we discuss the shortest representation of Hilbert's identity and its related tensor rank problem, in particular providing a lower bound for the number of terms in the identity. Finally, we conclude this paper with an application of determining the complexity of matrix $2 \mapsto 4$ norm problem, to illustrate the usefulness of our approach.
1.1. Notation. Throughout we adopt the notation of the lowercase letters to denote vectors (e.g., $x \in \mathbb{R}^{n}$ ), the capital letters to denote matrices (e.g., $A \in \mathbb{R}^{n^{2}}$ ), and the capital calligraphy letters to denote higher ( $\geq 3$ ) order tensors (e.g., $\mathscr{F} \in \mathbb{R}^{n^{4}}$ ), with subscriptions of indices being their entries (e.g., $x_{1}, A_{i j}, \mathscr{F}_{i_{1} i_{2} i_{3} i_{4}} \in \mathbb{R}$ ). A tensor is said to be supersymmetric if its entries are invariant under all permutations of its indices. As mentioned earlier, the symmetrization mapping "sym" makes a given tensor to be supersymmetric, which is $\mathscr{F}=\operatorname{sym}(\mathscr{G})$ with

$$
\mathscr{F}_{i_{1}, i_{2}, \ldots, i_{d}}=\frac{1}{\left|\Pi\left(i_{1}, i_{2}, \ldots, i_{d}\right)\right|} \sum_{\pi \in \Pi\left(i_{1}, i_{2}, \ldots, i_{d}\right)} \quad \mathcal{G}_{\pi} \quad \forall 1 \leq i_{1}, i_{2}, \ldots, i_{d} \leq n
$$

where $\Pi\left(i_{1}, i_{2}, \ldots, i_{d}\right)$ is the set of all distinct permutations of the indices $\left\{i_{1}, i_{2}, \ldots, i_{d}\right\}$. The symbol " $\otimes$ " represents the outer product of vectors or matrices. In particular, if

$$
\mathscr{F}=\underbrace{x \otimes x \otimes \cdots \otimes x}_{d} \quad \text { for some } x \in \mathbb{R}^{n},
$$

[^0]then $\mathscr{F}_{i_{1}, i_{2}, \ldots, i_{d}}=\prod_{k=1}^{d} x_{i_{k}}$; and if
$$
\mathscr{G}=\underbrace{X \otimes X \otimes \cdots \otimes X}_{d} \text { for some } X \in \mathbb{R}^{n^{2}},
$$
then $\mathscr{G}_{i_{1}, i_{2}, \ldots, i_{2 d}}=\prod_{k=1}^{d} X_{i_{2 k-}-i_{2} k}$. Besides, $\Delta$ denotes the supporting set of certain random variable, and $\Omega \subseteq \mathbb{R}^{n}$ is the sample space of a set of random variables $\left\{\xi_{1}, \xi_{2}, \ldots, \xi_{n}\right\}$, i.e., the space of all possible outcomes of $\left(\xi_{1}, \xi_{2}, \ldots, \xi_{n}\right)^{\mathrm{T}}$. Finally, the following two subsets of $\mathbb{Z}_{+}^{n}$ are frequently used in the discussion:
$$
\mathbb{P}_{k}^{n}:=\left\{\left(p_{1}, p_{2}, \ldots, p_{n}\right)^{\mathrm{T}} \in \mathbb{Z}_{+}^{n} \mid p_{1}+p_{2}+\cdots+p_{n}=k\right\},
$$
and for given prime number $q$,
$$
\mathbb{P}_{k}^{n}(q):=\left\{p \in \mathbb{P}_{k}^{n} \mid \exists i(1 \leq i \leq n) \text { such that } q \nmid p_{i}\right\} .
$$

It is easy to see that $\left|\mathbb{P}_{k}^{n}(q)\right| \leq\left|\mathbb{P}_{k}^{n}\right|=\binom{n+k-1}{k}$.
2. Hilbert's identity and $2 d$-th moments tensor. Let us start our discussion with the famous Hilbert's identity, which states that for any fixed positive integers $d$ and $n$, there always exist rational vectors $b^{1}, b^{2}, \ldots$, $b^{t} \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
\left(\sum_{i=1}^{n} x_{i}^{2}\right)^{d}=\sum_{j=1}^{t}\left(\left(b^{j}\right)^{\mathrm{T}} x\right)^{2 d} \quad \forall x=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{\mathrm{T}} \in \mathbb{R}^{n} . \tag{5}
\end{equation*}
$$

For instance, when $n=4$ and $d=2$, we have

$$
\begin{equation*}
\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}+x_{4}^{2}\right)^{2}=\frac{1}{6} \sum_{1 \leq i<j \leq 4}\left(x_{i}+x_{j}\right)^{4}+\frac{1}{6} \sum_{1 \leq i<j \leq 4}\left(x_{i}-x_{j}\right)^{4}, \tag{6}
\end{equation*}
$$

which is called Liouville's identity. It is worth mentioning that Hilbert's identity is very well known and is a fundamental result in mathematics. For example, with the help of (5), Reznick [18] managed to prove the result: Let $p(x)$ be $2 d$-th degree homogeneous positive polynomial in $x \in \mathbb{R}^{n}$. Then there exists a positive integer $r$ and vectors $b^{1}, b^{2}, \ldots, b^{r} \in \mathbb{R}^{n}$ such that

$$
\|x\|_{2}^{2 r-2 d} p(x)=\sum_{i=1}^{r}\left(\left(b^{i}\right)^{\mathrm{T}} x\right)^{2 r} .
$$

Reznick's result solved Hilbert's 17th problem constructively (albeit only for the case $p(x)$ being positive definite). As another example, Hilbert [10] in 1909 solved Waring's problem in the affirmative for all $k$ : Can every positive integer be expressed as a sum of at most $g(k) k$-th powers of positive integers, where $g(k)$ depends only on $k$, not on the number being represented? The key underpinning tool in the proof is also Hilbert's identity (5); see, e.g., Ellison [7] and Pollack [16] for more stories on Warning's problem and Hilbert's identity. In fact, Hilbert's identity can be readily extended to a more general setting. For any given $A \succeq 0$, by letting $y=A^{1 / 2} x$ and applying (5), one has

$$
\left(x^{\mathrm{T}} A x\right)^{2}=\left(y^{\mathrm{T}} y\right)^{2}=\sum_{j=1}^{t}\left(\left(b^{j}\right)^{\mathrm{T}} y\right)^{2 d}=\sum_{j=1}^{t}\left(\left(b^{j}\right)^{\mathrm{T}} A^{1 / 2} x\right)^{2 d},
$$

which guarantees the existence of vectors $a^{1}, a^{2}, \ldots, a^{t} \in \mathbb{R}^{n}$ with $a^{j}=A^{1 / 2} b^{j}$ for $j=1,2, \ldots, t$ such that

$$
\begin{equation*}
\left(x^{\mathrm{T}} A x\right)^{d}=\sum_{j=1}^{t}\left(\left(a^{j}\right)^{\mathrm{T}} x\right)^{2 d} . \tag{7}
\end{equation*}
$$

The discussion so far appears to be only concerned about decomposing a specific polynomial function. Let us now relate Hilbert's identity to the moments tensor. Observe that supersymmetric tensors are bijectively related to homogenous polynomial functions. In particular, if

$$
f(x)=\sum_{1 \leq i_{1} \leq i_{2} \leq \cdots \leq i_{d} \leq n} \mathscr{G}_{i_{1}, i_{2}, \ldots, i_{d}} \prod_{k=1}^{d} x_{i_{k}}
$$

is a $d$-th degree homogenous polynomial, then its associated supersymmetric tensor $\mathscr{F}$ with

$$
\mathscr{F}_{i_{1}, i_{2}, \ldots, i_{d}}=\frac{\mathscr{G}_{i_{1}, i_{2}, \ldots, i_{d}}}{\left|\Pi\left(i_{1}, i_{2}, \ldots, i_{d}\right)\right|}
$$

is uniquely determined by

$$
f(x)=\mathscr{F}(\underbrace{x, x, \ldots, x}_{d}),
$$

and vice versa. This is the same as the one-to-one correspondence between symmetric matrices and quadratic forms. Therefore, the tensor

$$
\operatorname{sym}(\underbrace{A \otimes A \otimes \cdots \otimes A}_{d})
$$

is associated with the polynomial $\left(x^{\mathrm{T}} A x\right)^{d}$, and the following relationship holds immediately.
Proposition 2.1. For any $A \succeq 0$, there exist vectors $a^{1}, a^{2}, \ldots, a^{t} \in \mathbb{R}^{n}$ such that

$$
\left(x^{\mathrm{T}} A x\right)^{2}=\sum_{j=1}^{t}\left(\left(a^{j}\right)^{\mathrm{T}} x\right)^{2 d}
$$

i.e.,

$$
\operatorname{sym}(\underbrace{A \otimes A \otimes \cdots \otimes A}_{d})=\sum_{i=1}^{t} \underbrace{a^{i} \otimes a^{i} \otimes \cdots \otimes a^{i}}_{2 d}
$$

This implies that tensor $\operatorname{sym}(\underbrace{A \otimes A \otimes \cdots \otimes A}_{d})$ is a $2 d$-th moments tensor if $A \succeq 0$.
As we mentioned earlier, the size of such representation from Hilbert's identity is exponential in $n$. To see this, let us recall the claim of Hilbert (see Nathanson [14]): Given fixed positive integers $d$ and $n$, there exist $2 d+1$ real numbers $\beta_{1}, \beta_{2}, \ldots, \beta_{2 d+1}, 2 d+1$ positive real numbers $\rho_{1}, \rho_{2}, \ldots, \rho_{2 d+1}$, and a positive real number $\alpha_{d}$, such that

$$
\begin{equation*}
\left(x^{\mathrm{T}} x\right)^{d}=\frac{1}{\alpha_{d}} \sum_{i_{1}=1}^{n} \sum_{i_{2}=1}^{n} \cdots \sum_{i_{2 d+1}=1}^{n} \rho_{i_{1}}, \rho_{i_{2}}, \ldots, \rho_{i_{2 d+1}}\left(\beta_{i_{1}} x_{1}+\beta_{i_{2}} x_{2}+\cdots+\beta_{i_{2 d+1}} x_{i_{2 d+1}}\right)^{2 d} \tag{8}
\end{equation*}
$$

It is obvious that the number of $2 d$-powered linear terms on the right-hand side of $(8)$ is $(2 d+1)^{n}$, which is too lengthy for practical purposes. In the following, let us focus on how to get a polynomial-size decomposition of Hilbert's identity, or essentially the tensor $\operatorname{sym}(\underbrace{A \otimes A \otimes \cdots \otimes A})$ with $A \succeq 0$.

In light of the above discussion, it suffices to find a polynomial-size representation of (5). Toward this end, let us first rewrite $\left(x^{\mathrm{T}} x\right)^{d}$ in terms of the expectation of a polynomial function. In particular, by defining i.i.d. random variables $\xi_{1}, \xi_{2}, \ldots, \xi_{n}$ with supporting set $\Delta=\left\{\beta_{1}, \beta_{2}, \ldots, \beta_{2 d+1}\right\}$ and $\operatorname{Prob}\left(\xi_{k}=\beta_{i}\right)=\rho_{i} / \gamma_{d}$ for all $1 \leq i \leq 2 d+1$ and $1 \leq k \leq n$, where $\gamma_{d}=\sum_{i=1}^{2 d+1} \beta_{i}$, identity (8) is equivalent to

$$
\begin{equation*}
\left(x^{\mathrm{T}} x\right)^{d}=\frac{\gamma_{d}^{d}}{\alpha_{d}} \mathrm{E}\left[\left(\sum_{j=1}^{n} \xi_{j} x_{j}\right)^{2 d}\right]=\frac{\gamma_{d}^{d}}{\alpha_{d}} \sum_{p \in \mathbb{P}_{2 d}^{n}} \mathrm{E}\left[\prod_{j=1}^{n} \xi_{j}^{p_{j}}\right] \prod_{j=1}^{n} x_{j}^{p_{j}}=\frac{\gamma_{d}^{d}}{\alpha_{d}} \sum_{p \in \mathbb{P}_{2 d}{ }^{n}} \prod_{j=1}^{n} \mathrm{E}\left[\xi_{j}^{p_{j}}\right] \prod_{j=1}^{n} x_{j}^{p_{j}} . \tag{9}
\end{equation*}
$$

As a consequence, if for any $n$ random variables $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ satisfying

$$
\begin{equation*}
\mathrm{E}\left[\prod_{j=1}^{n} \eta_{j}^{p_{j}}\right]=\prod_{j=1}^{n} \mathrm{E}\left[\eta_{j}^{p_{j}}\right] \quad \forall p \in \mathbb{P}_{2 d}^{n} \tag{10}
\end{equation*}
$$

and $\mathrm{E}\left[\eta_{j}^{p}\right]=\mathrm{E}\left[\xi_{1}^{p}\right]$ for all $0<p \leq 2 d$ and $1 \leq j \leq n$, then it is straightforward to verify that $\left(x^{\mathrm{T}} x\right)^{d}=$ $\left(\gamma_{d}^{d} / \alpha_{d}\right) \mathrm{E}\left[\left(\sum_{j=1}^{n} \eta_{j} x_{j}\right)^{2 d}\right]$. Notice that (10) is actually equivalent to $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ being $2 d$-wise uncorrelated, and we have the next result following (9) and (10).

Proposition 2.2. If $\xi_{1}, \xi_{2}, \ldots, \xi_{n}$ are i.i.d. random variables, and $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ are $2 d$-wise uncorrelated, satisfying the moments constraints $\mathrm{E}\left[\eta_{j}^{p}\right]=\mathrm{E}\left[\xi_{1}^{p}\right]$ for all $0<p \leq 2 d$ and $1 \leq j \leq n$, then $\mathrm{E}\left[\left(\sum_{j=1}^{n} \xi_{j} x_{j}\right)^{2 d}\right]=$ $\mathrm{E}\left[\left(\sum_{j=1}^{n} \eta_{j} x_{j}\right)^{2 d}\right]$.

We end this section with the conclusion that the key to reducing the length of representation in (5) is to construct $2 d$-wise uncorrelated random variables satisfying certain moments conditions, such that the sample space is as small as possible, which will be the subject of our subsequent discussions. As we will see later, the construction makes use of the structure of the support set (4). For general support sets, the techniques considered in Karloff and Mansour [13] may be useful, and it is a topic for future research.
3. Construction of $k$-wise uncorrelated random variables. In this section, we shall construct $k$-wise uncorrelated random variables, which are identical and uniformly distributed on $\Delta_{q}$ defined by (4). The rough idea is as follows. We first generate $m$ i.i.d. random variables $\xi_{1}, \xi_{2}, \ldots, \xi_{m}$, based on which we can define new random variables $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ such that $\eta_{i}:=\prod_{1 \leq j \leq m} \xi_{j}^{c_{i j}}$ for $i=1,2, \ldots, n$. Therefore, the size of the sample space of $\left\{\eta_{1}, \eta_{2}, \ldots, \eta_{n}\right\}$ is bounded above by $q^{m}$, which yields a polynomial-size space if we let $m=O\left(\log _{q} n\right)$. The remaining part of this section is devoted to the discussion of the property for the power indices $c_{i j}$ 's, in order to guarantee $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ to be $k$-wise uncorrelated, and how to find those power indices.
3.1. $k$-wise regular sequence. Let us start with some notations and definitions for the preparation. Suppose $c$ is a number with $m$ digits and $c[l]$ is the value of its $l$-th bit. We call $c$ to be endowed with the base $q$, if $c[l] \in\{0,1, \ldots, q-1\}$ for all $1 \leq l \leq m$. In other words, $c=\sum_{l=1}^{m} c[l] q^{l-1}$. Now we can define the concept of $k$-wise regular sequence as follows.

Definition 3.1. A sequence of $m$ digits numbers $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ of base $q$ is called $k$-wise regular if for any $p \in \mathbb{P}_{k}^{n}(q)$, there exists $l(1 \leq l \leq m)$ such that

$$
\sum_{j=1}^{n} p_{j} \cdot c_{j}[l] \neq 0 \bmod q
$$

Why are we interested in such regular sequences? The answer lies in the following proposition.
Proposition 3.1. Suppose $m$ digits numbers $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ of base $q$ are $k$-wise regular, where $q$ is a prime number, and $\xi_{1}, \xi_{2}, \ldots, \xi_{m}$ are i.i.d. random variables uniformly distributed on $\Delta_{q}$. Then $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ with

$$
\begin{equation*}
\eta_{i}:=\prod_{l=1}^{m} \xi_{l}^{c_{i}[l]} \quad \text { for } i=1,2, \ldots, n \tag{11}
\end{equation*}
$$

are $k$-wise uncorrelated.
Proof. Let $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ be defined as in (11). As $\xi_{i}$ is uniformly distributed on $\Delta_{q}$ for $1 \leq i \leq m$ and $q$ is prime, we have

$$
\mathrm{E}\left[\xi_{i}^{p}\right]=\mathrm{E}\left[\eta_{j}^{p}\right]= \begin{cases}1 & \text { if } q \mid p \\ 0 & \text { otherwise }\end{cases}
$$

for any $i$ and any $j$ with $c_{j} \neq(0,0, \ldots, 0)$. Otherwise if $c_{j}=(0,0, \ldots, 0)$ for some $j$, then $\mathrm{E}\left[\eta_{j}^{p}\right]=1$.
For any given $p \in \mathbb{P}_{k}^{n}$, if $q \mid p_{i}$ for all $1 \leq i \leq n$, then

$$
\mathrm{E}\left[\prod_{j=1}^{n} \eta_{j}^{p_{j}}\right]=\mathrm{E}\left[\left(\prod_{l=1}^{m} \xi_{l}^{p_{1} \cdot c_{1}[l]}\right)\left(\prod_{l=1}^{m} \xi_{l}^{p_{2} \cdot c_{2}[l]}\right) \cdots\left(\prod_{l=1}^{m} \xi_{l}^{p_{n} \cdot c_{n}[l]}\right)\right]=\prod_{l=1}^{m} \mathrm{E}\left[\xi_{l}^{\sum_{j=1}^{n} p_{j} \cdot c_{j}[l]}\right]=1=\prod_{j=1}^{n} \mathrm{E}\left[\eta_{j}^{p_{j}}\right] .
$$

Otherwise, there exists some $i_{0}$ such that $q \nmid p_{i_{0}}$, implying that $p \in \mathbb{P}_{k}^{n}(q)$. By $k$-wise regularity, we can find some $l_{0}$ satisfying $\sum_{j=1}^{n} p_{j} \cdot c_{j}\left[l_{0}\right] \neq 0 \bmod q$, implies that $\mathrm{E}\left[\xi_{l_{0}}^{\sum_{j=1}^{n} p_{j} \cdot c_{j}\left[l_{0}\right]}\right]=0$. Moreover, there exists some $j_{0}$ such that $p_{j_{0}} \cdot c_{j_{0}}\left[l_{0}\right] \neq 0 \bmod q$, i.e., $q \nmid p_{j_{0}}$ and $c_{j_{0}}\left[l_{0}\right] \neq 0$. This leads to $\mathrm{E}\left[\eta_{j_{0}}^{p_{j_{0}}}\right]=0$, and we have

$$
\mathrm{E}\left[\prod_{j=1}^{n} \eta_{j}^{p_{j}}\right]=\prod_{l=1}^{m} \mathrm{E}\left[\xi_{l}^{\sum_{j=1}^{n} p_{j} \cdot c_{j}[l]}\right]=0=\prod_{j=1}^{n} \mathrm{E}\left[\eta_{j}^{p_{j}}\right]
$$

and the conclusion follows.
3.2. A randomized algorithm. We shall now focus on how to find such $k$-wise regular sequence $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ of base $q$. First, we present a randomized process, in which $c_{i}[l]$ is randomly and uniformly chosen from $\{0,1, \ldots, q-1\}$ for all $1 \leq i \leq n$ and $1 \leq l \leq m$. The algorithm is as follows.

## Algorithm RAN <br> Input: Dimension $n$ and $m:=\left\lceil k \log _{q} n\right\rceil$.

Output: A sequence $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ in $m$ digits of base $q$.
Step 0. Construct $S=\{(\underbrace{0, \ldots, 0,0}_{m}),(\underbrace{0, \ldots, 0,1}_{m}), \ldots,(\underbrace{q-1, \ldots, q-1, q-1}_{m})\}$ of base $q$.
Step 1. Independently and uniformly take $c_{i} \in S$ for $i=1,2, \ldots, n$.
Step 2. Assemble the sequence $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ and exit.

Theorem 3.1. If $1<k<n$ and $q$ is a prime number, then Algorithm RAN returns a $k$-wise m-digit regular sequence $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ of base $q$ with probability at least $1-(1.5)^{k-1} / k!$, which is independent of $n$ and $q$.

Proof. Since $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ is a sequence of $m$-digit numbers of base $q$, if it is not regular, then there exists $p \in \mathbb{P}_{k}^{n}$, such that

$$
\sum_{j=1}^{n} p_{j} \cdot c_{j}[l]=0 \bmod q \quad \forall 1 \leq l \leq m
$$

Therefore, we have

$$
\operatorname{Prob}\left\{\left\{c_{1}, c_{2}, \ldots, c_{n}\right\} \text { is not } k \text {-wise regular }\right\} \leq \sum_{p \in \mathbb{P}_{k}^{n}(q)} \operatorname{Prob}\left\{\sum_{j=1}^{n} p_{j} \cdot c_{j}[l]=0 \bmod q, \forall 1 \leq l \leq m\right\}
$$

For any given $p \in \mathbb{P}_{k}^{n}(q)$, we may without loss of generality assume that $q \nmid p_{n}$. If we fix $c_{1}, c_{2}, \ldots, c_{n-1}$, as $q$ is prime, then there is only one solution for $c_{n}$ such that $\sum_{j=1}^{n} p_{j} \cdot c_{j}[l]=0 \bmod q$ for all $1 \leq l \leq m$. Combining the fact that $c_{1}, c_{2}, \ldots, c_{n}$ are independently and uniformly generated, we have

$$
\begin{align*}
& \operatorname{Prob}\left\{\sum_{j=1}^{n} p_{j} \cdot c_{j}[l]=0 \bmod q, \forall 1 \leq l \leq m\right\} \\
& \quad=\operatorname{Prob}\left\{\sum_{j=1}^{n} p_{j} \cdot c_{j}[l]=0 \bmod q, \forall 1 \leq l \leq m \mid c_{1}=d_{1}, c_{2}=d_{2}, \ldots, c_{n-1}=d_{n-1}\right\} \\
& \quad \cdot \sum_{d_{1}, d_{2}, \ldots, d_{n-1} \in S} \operatorname{Prob}\left\{c_{1}=d_{1}, c_{2}=d_{2}, \ldots, c_{n-1}=d_{n-1}\right\} \\
& \quad=\frac{1}{q^{m}} \sum_{d_{1}, d_{2}, \ldots, d_{n-1} \in S} \operatorname{Prob}\left\{c_{1}=d_{1}, c_{2}=d_{2}, \ldots, c_{n-1}=d_{n-1}\right\} \\
& \quad \leq \frac{1}{n^{k}} \tag{12}
\end{align*}
$$

Finally,

$$
\begin{aligned}
& \operatorname{Prob}\left\{\left\{c_{1}, c_{2}, \ldots, c_{n}\right\} \text { is } k \text {-wise regular }\right\} \\
& \quad=1-\operatorname{Prob}\left\{\left\{c_{1}, c_{2}, \ldots, c_{n}\right\} \text { is not } k \text {-wise regular }\right\} \\
& \quad \geq 1-\left|\mathbb{P}_{k}^{n}(q)\right| \cdot \frac{1}{n^{k}} \geq 1-\left|\mathbb{P}_{k}^{n}\right| \cdot \frac{1}{n^{k}}=1-\binom{n+k-1}{k} \cdot \frac{1}{n^{k}} \geq 1-\frac{(1.5)^{k-1}}{k!}
\end{aligned}
$$

For some special $q$ and $k$, in particular relating to the simplest case of Hilbert's identity (4-wise regular sequence of base 2), the lower bound of the probability in Theorem 3.1 can be improved.

Proposition 3.2. If $k=4$ and $q=2$, then Algorithm RAN returns a 4 -wise regular sequence $\left\{c_{1}, c_{2}\right.$, $\left.\ldots, c_{n}\right\}$ of base 2 with probability at least $1-1 /\left(2 n^{2}\right)-1 / 4$ !.

The proof is similar to that of Theorem 3.1, and thus is omitted.
3.3. Derandomization. Although $k$-wise regular sequence always exists and can be found with high probability, one may however wish to construct such regular sequence deterministically. In fact, this is possible if we apply Theorem 3.1 in a slightly different manner, which is shown in the following algorithm. Basically, we start with a short regular sequence $C$, and enumerate all the remaining numbers in order to find $c$ such that $C \cup\{c\}$ is also regular. Updating $C$ with $C \cup\{c\}$, we repeat this procedure until the cardinality of $C$ reaches $n$. Moreover, thanks to the polynomial-size sample space, this "brute force" approach still runs in polynomial-time.

## Algorithm DET

Input: Dimension $n$ and $m:=\left\lceil k \log _{q} n\right\rceil$.
Output: A sequence $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ in $m$ digits of base $q$.
Step 0 . Construct $S=\{(\underbrace{0, \ldots, 0,0}_{m}),(\underbrace{0, \ldots, 0,1}_{m}), \ldots,(\underbrace{q-1, \ldots, q-1, q-1}_{m})\}$ of base $q$, and a sequence $C:=\left\{c_{1}, c_{2}, \ldots, c_{k}\right\}$ in $m$ digits, where $c_{i}:=(0, \ldots, 0,0,1, \underbrace{0, \ldots, 0,0}_{k-1})$ for
$i=1,2, \ldots, k$. Let the index count be $\tau:=k$.
Step 1. If $\tau=n$, then go to Step 2; Otherwise enumerate $S \backslash C$ to find a $c \in S \backslash C$ such that $C \cup\{c\}$ is $k$-wise regular. Let $c_{\tau+1}:=c, C:=C \cup\left\{c_{\tau+1}\right\}$ and $\tau:=\tau+1$, and return to Step 1.
Step 2. Assemble the sequence $\left\{c_{1}, c_{2}, \ldots, c_{n}\right\}$ and exit.
It is obvious that the initial sequence $\left\{c_{1}, c_{2}, \ldots, c_{k}\right\}$ is $k$-wise regular. For Algorithm $D E T$ to terminate successfully, it remains to argue that it is always possible to expand the $k$-wise regular sequence by one in Step 1 , as long as $\tau<n$.

Theorem 3.2. Suppose that $3 \leq k \leq \tau<n, q$ is a prime number, and $C$ with $|C|=\tau$ is $k$-wise regular. If we uniformly pick $c_{\tau+1}$ from $S$, then

$$
\operatorname{Prob}\left\{C \cup\left\{c_{\tau+1}\right\} \text { is } k \text {-wise regular }\right\} \geq 1-\frac{(1.5)^{k}}{k!}\left(\frac{\tau+1}{n}\right)^{k}
$$

ensuring that $\left\{c_{\tau+1} \in S \mid C \cup\left\{c_{\tau+1}\right\}\right.$ is $k$-wise regular $\} \neq \varnothing$.
Proof. Like in the proof of Theorem 3.1, we have

$$
\operatorname{Prob}\left\{C \cup\left\{c_{\tau+1}\right\} \text { is not } k \text {-wise regular }\right\} \leq \sum_{p \in \mathbb{P}_{k}^{\tau+1}(q)} \operatorname{Prob}\left\{\sum_{j=1}^{\tau+1} p_{j} \cdot c_{j}[l]=0 \bmod q, \forall 1 \leq l \leq m\right\}
$$

For any $p \in \mathbb{P}_{k}^{\tau+1}(q)$, since $q$ is prime, by using a similar argument as of (12), we can get

$$
\operatorname{Prob}\left\{\sum_{j=1}^{\tau+1} p_{j} \cdot c_{j}[l]=0 \bmod q, \forall 1 \leq l \leq m\right\} \leq \frac{1}{n^{k}}
$$

Essentially, the argument in (12) works by conditioning on the elements in $C$, the selection ordering in $C$ during the previous steps is not important. Therefore,

$$
\operatorname{Prob}\left\{C \cup\left\{c_{\tau+1}\right\} \text { is } k \text {-wise regular }\right\} \geq 1-\left|\mathbb{P}_{k}^{\tau+1}(q)\right| \frac{1}{n^{k}} \geq 1-\binom{\tau+k}{k} \frac{1}{n^{k}} \geq 1-\frac{(1.5)^{k}}{k!}\left(\frac{\tau+1}{n}\right)^{k}>0
$$

By the above theorem, Step 1 of Algorithm DET guarantees to expand the $k$-wise regular sequence of base $q$ before reaching the desired cardinality $\tau=n$. A straightforward computation shows that Algorithm DET requires an overall complexity of $O\left(n^{2 k-1} \log _{q} n\right)$.

## 4. Polynomial-size representation of moments tensor

4.1. Polynomial-size representation of the fourth moments tensor. With the help of $k$-wise uncorrelated random variables, we are able to construct polynomial-size representation of the fourth moments tensor. In Hilbert's construction (9), the support set $\Delta$ is too general to apply the result in $\S 3$. However as we mentioned earlier, such decomposition of (9) is not unique. In fact, when $d=2$, we observe that

$$
\begin{equation*}
\left(x^{\mathrm{T}} x\right)^{2}=\left(\sum_{i=1}^{n} x_{i}^{2}\right)^{2}=\frac{2}{3} \sum_{i=1}^{n} x_{i}^{4}+\frac{1}{3} \mathrm{E}\left[\left(\sum_{j=1}^{n} \xi_{j} x_{j}\right)^{4}\right], \tag{13}
\end{equation*}
$$

where $\xi_{1}, \xi_{2}, \ldots, \xi_{n}$ are i.i.d. symmetric Bernoulli random variables. Applying either Algorithm RAN or Algorithm DET leads to a 4-wise regular sequence of base 2, based on which we can define random variables $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ as we did in (11). Proposition 3.1 guarantees that $\eta_{1}, \eta_{2}, \ldots, \eta_{n}$ are 4-wise uncorrelated, and it is easy to check that

$$
\mathrm{E}\left[\eta_{j}\right]=\mathrm{E}\left[\eta_{j}^{3}\right]=\mathrm{E}\left[\xi_{1}\right]=\mathrm{E}\left[\xi_{1}^{3}\right]=0, \quad \mathrm{E}\left[\eta_{j}^{2}\right]=\mathrm{E}\left[\eta_{j}^{4}\right]=\mathrm{E}\left[\dot{\xi}_{1}^{2}\right]=\mathrm{E}\left[\xi_{1}^{4}\right]=1 \quad \forall 1 \leq j \leq n
$$

Thus, by Proposition 2.2, we have $\mathrm{E}\left[\left(\sum_{j=1}^{n} \eta_{j} x_{j}\right)^{4}\right]=\mathrm{E}\left[\left(\sum_{j=1}^{n} \xi_{j} x_{j}\right)^{4}\right]$. Moreover, the size of the sample space of $\left\{\eta_{1}, \eta_{2}, \ldots, \eta_{n}\right\}$ is at most $2^{\left[k \log _{q} n\right\rceil} \leq 2 n^{4}$, which means the new representation has at most $n+2 n^{4}$ fourth powered terms. Combining with Proposition 2.1, we have the following main result.

Theorem 4.1. Given a positive integer $n$, we can find $\tau\left(\leq 2 n^{4}\right)$ vectors $b^{1}, b^{2}, \ldots, b^{\tau} \in \mathbb{R}^{n}$ in polynomial time, such that

$$
\left(x^{\mathrm{T}} x\right)^{2}=\frac{2}{3} \sum_{i=1}^{n} x_{i}^{4}+\sum_{j=1}^{\tau}\left(\left(b^{j}\right)^{\mathrm{T}} x\right)^{4} \quad \forall x \in \mathbb{R}^{n},
$$

or equivalently,

$$
\operatorname{sym}(I \otimes I)=\frac{2}{3} \sum_{i=1}^{n} e^{i} \otimes e^{i} \otimes e^{i} \otimes e^{i}+\sum_{j=1}^{\tau} b^{j} \otimes b^{j} \otimes b^{j} \otimes b^{j}
$$

where $e^{i} \in \mathbb{R}^{n}$ is the $i$-th unit vector (with the $i$-th entry 1 and other entries zeros).
The result can be extended to a more general setting as follows.
Corollary 4.1. Given a positive semidefinite matrix $A \in \mathbb{R}^{n \times n}$, we can find $\tau\left(\leq 2 n^{4}+n\right)$ vectors $a^{1}, a^{2}, \ldots, a^{\tau} \in \mathbb{R}^{n}$ in polynomial time, such that

$$
\left(x^{\mathrm{T}} A x\right)^{2}=\sum_{i=1}^{\tau}\left(\left(a^{i}\right)^{\mathrm{T}} x\right)^{4} \quad \forall x \in \mathbb{R}^{n},
$$

or equivalently,

$$
\operatorname{sym}(A \otimes A)=\sum_{i=1}^{\tau} a^{i} \otimes a^{i} \otimes a^{i} \otimes a^{i}
$$

Proof. Because of the one to one correspondence between supersymmetric tensors and homogeneous polynomials, we only need to prove the first identity. By letting $y=A^{1 / 2} x$ and applying Theorem 4.1, we can find $b^{1}, b^{2}, \ldots, b^{\tau}$ in polynomial time with $\tau \leq 2 n^{4}$, such that

$$
\left(x^{\mathrm{T}} A x\right)^{2}=\left(y^{\mathrm{T}} y\right)^{2}=\frac{2}{3} \sum_{i=1}^{n} y_{i}^{4}+\sum_{j=1}^{\tau}\left(\left(b^{j}\right)^{\mathrm{T}} y\right)^{4}=\sum_{i=1}^{n}\left(\left(\frac{2}{3}\right)^{1 / 4}\left(e^{i}\right)^{\mathrm{T}} A^{1 / 2} x\right)^{4}+\sum_{j=1}^{\tau}\left(\left(b^{j}\right)^{\mathrm{T}} A^{1 / 2} x\right)^{4}
$$

The conclusion follows by letting $a^{i}=\left(\frac{2}{3}\right)^{1 / 4} A^{1 / 2} e^{i}$ for $i=1,2, \ldots, n$, and $a^{i+n}=A^{1 / 2} b^{i}$ for $i=1,2, \ldots, \tau$.
4.2. Polynomial-size representation of complex $q d$-th moments tensor. In this subsection we are going to generalize the result in $\S 4.1$ to $q d$-th moments tensor. Denote $\mathscr{J}^{q}$ to be the $q$-th order identity tensor, whose entry is 1 when all its indices are identical, and is zero otherwise. We are interested in whether

$$
\underbrace{\mathscr{I}^{q} \otimes \mathscr{J}^{q} \otimes \cdots \otimes \mathscr{J}^{q}}_{d}
$$

is a $q d$-th moments tensor or not. If it is true, then for any given positive integers $q, d$ and $n$, there exist vectors $a^{1}, a^{2}, \ldots, a^{t} \in \mathbb{R}^{n}$, such that

$$
\begin{equation*}
\operatorname{sym}(\underbrace{\mathscr{F}^{q} \otimes \mathscr{J}^{q} \otimes \cdots \otimes \mathscr{J}^{q}}_{d})=\sum_{i=1}^{t} \underbrace{a^{i} \otimes a^{i} \otimes \cdots \otimes a^{i}}_{q d}, \tag{14}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\left(\sum_{i=1}^{n} x_{i}^{q}\right)^{d}=\sum_{j=1}^{t}\left(\left(a^{j}\right)^{\mathrm{T}} x\right)^{q d} \quad \forall x \in \mathbb{R}^{n} . \tag{15}
\end{equation*}
$$

Unfortunately, the above does not hold in general, as the following counterexample shows.
Example 4.1. The function $f(x)=\left(x_{1}^{3}+x_{2}^{3}\right)^{2}=x_{1}^{6}+2 x_{1}^{3} x_{2}^{3}+x_{2}^{6}$ cannot be decomposed in the form of (15) with $q=3$ and $d=2$, i.e., a sum of sixth powered linear terms.

This can be easily proven by contradiction. Suppose we can find $a^{1}, a^{2}, \ldots, a^{t} \in \mathbb{R}^{n}$, such that

$$
\begin{equation*}
x_{1}^{6}+2 x_{1}^{3} x_{2}^{3}+x_{2}^{6}=\sum_{i=1}^{t}\left(a_{i} x_{1}+b_{i} x_{2}\right)^{6} . \tag{16}
\end{equation*}
$$

There must exist some $\left(a_{j}, b_{j}\right)$ with $a_{j} b_{j} \neq 0$, since otherwise there is no monomial $x_{1}^{3} x_{2}^{3}$ in the right-hand side of (16). As a consequence, the coefficient of monomial $x_{1}^{2} x_{2}^{4}$ in the right-hand side of (16) is at least $\binom{6}{2} a_{j}^{2} b_{j}^{4}>0$, which is null on the left side of the equation, leading to a contradiction.

In the same vein one can actually show that (15) cannot hold for any $q \geq 3$. Therefore, we turn to $q d$-th moments tensor in the complex domain, i.e., both entries of the tensor and vector $a^{i}$ 's in (14) and (15) are now allowed to take complex values. Similar to (13), we have the following identity:

$$
\begin{equation*}
\left(\sum_{j=1}^{n} x_{j}^{q}\right)^{2}=\left(1-\frac{2}{\binom{2 q}{q}}\right) \sum_{j=1}^{n} x_{j}^{2 q}+\frac{2}{\binom{2 q}{q}} \mathrm{E}\left[\left(\sum_{i=1}^{n} \xi_{i} x_{i}\right)^{2 q}\right], \tag{17}
\end{equation*}
$$

where $\xi_{1}, \xi_{2}, \ldots, \xi_{n}$ are i.i.d. random variables uniformly distributed on $\Delta_{q}$. Moreover, we can further prove (15) for the more general complex case.

Proposition 4.1. For any given positive integers $q$, $d$, and $n$, there exist $a^{1}, a^{2}, \ldots, a^{\tau} \in \mathbb{C}^{n}$ such that

$$
\begin{equation*}
\left(\sum_{i=1}^{n} x_{i}^{q}\right)^{2^{d}}=\sum_{j=1}^{\tau}\left(\left(a^{j}\right)^{\mathrm{T}} x\right)^{2^{d} q} \quad \forall x \in \mathbb{C}^{n}, \tag{18}
\end{equation*}
$$

or equivalently,

$$
\operatorname{sym}(\underbrace{\mathscr{S}^{q} \otimes \mathscr{J}^{q} \otimes \cdots \otimes \mathscr{J}^{q}}_{2^{d}})=\sum_{i=1}^{t} \underbrace{a^{i} \otimes a^{i} \otimes \cdots \otimes a^{i}}_{2^{d} q} .
$$

Proof. Because of the one to one correspondence between supersymmetric tensors and homogeneous polynomials, we only need to prove the first identity, whose proof is based on mathematical induction. The case $d=1$ is already guaranteed by (17). Suppose that (18) is true for $d-1$, then there exist $b^{1}, b^{2}, \ldots, b^{t} \in \mathbb{C}^{n}$ such that

$$
\left(\sum_{i=1}^{n} x_{i}^{q}\right)^{2^{d}}=\left(\left(\sum_{i=1}^{n} x_{i}^{q}\right)^{2^{d-1}}\right)^{2}=\left(\sum_{j=1}^{t}\left(\left(b^{j}\right)^{\mathrm{T}} x\right)^{2^{d-1} q}\right)^{2} .
$$

By applying (17) to the above identity, there exist $c^{1}, c^{2}, \ldots, c^{\tau} \in \mathbb{C}^{t}$, such that

$$
\left(\sum_{i=1}^{n} x_{i}^{q}\right)^{2^{d}}=\left(\sum_{j=1}^{t}\left(\left(b^{j}\right)^{\mathrm{T}} x\right)^{2^{d-1} q}\right)^{2}=\sum_{i=1}^{\tau}\left(\sum_{j=1}^{t}\left(c^{i}\right)_{j} \cdot\left(b^{j}\right)^{\mathrm{T}} x\right)^{2^{d} q}=\sum_{i=1}^{\tau}\left(\left(c^{i}\right)^{\mathrm{T}} B^{\mathrm{T}} x\right)^{2^{d} q},
$$

where $B=\left(b^{1}, b^{2}, \ldots, b^{t}\right) \in \mathbb{C}^{n \times t}$. Letting $a^{i}=B c^{i}(1 \leq i \leq \tau)$ completes the inductive step.
The next step is to reduce the number $\tau$ in (18). Under the condition that $q$ is prime, we can get a $k$-wise regular sequence of base $q$ using either Algorithm RAN or Algorithm DET. With the help of Theorem 2.2, we can further get a polynomial-size representation of complex Hilbert's identity and complex $2^{d} q$-th moments tensor, by applying a similar argument as in Theorem 4.1.

Theorem 4.2. For any given positive integers $q$, $d$, and $n$ with $q$ being prime, we can find $\tau \leq O\left(n^{(2 q)^{2 d-1}}\right)$ vectors $a^{1}, a^{2}, \ldots, a^{\tau} \in \mathbb{C}^{n}$ in polynomial time, such that

$$
\left(\sum_{i=1}^{n} x_{i}^{q}\right)^{2^{d}}=\sum_{i=1}^{\tau}\left(\left(a^{i}\right)^{\mathrm{T}} x\right)^{2^{d} q} \quad \forall x \in \mathbb{C}^{n},
$$

or equivalently,

$$
\operatorname{sym}(\underbrace{\mathscr{I}^{q} \otimes \mathscr{J}^{q} \otimes \cdots \otimes \mathscr{J}^{q}}_{2^{d}})=\sum_{i=1}^{\tau} \underbrace{a^{i} \otimes a^{i} \otimes \cdots \otimes a^{i}}_{2^{d} q}
$$

5. Shortest representation of Hilbert's identity. In §4.1, we constructed polynomial-size representation of Hilbert's identity, in particular, the fourth moments tensor sym $(I \times I)$. The number of fourth powered linear functions required (in Theorem 4.1) is $n+2 n^{4}$. As we shall see later, this size is in general not the smallest possible. This raises the issue of how to find the shortest representation of the fourth moments tensor. In general, we are interested in the following quantity:

$$
\tau_{2 d}(n):=\min _{m \in \mathbb{Z}_{+}}\left\{\exists b^{1}, b^{2}, \ldots, b^{m} \in \mathbb{R}^{n}, \text { such that }\left(x^{\mathrm{T}} x\right)^{d}=\sum_{i=1}^{m}\left(\left(b^{i}\right)^{\mathrm{T}} x\right)^{2 d} \forall x \in \mathbb{R}^{n}\right\} .
$$

If fact, $\tau_{2 d}(n)$ is closely related to the rank of the supersymmetric tensor $\operatorname{sym}(\underbrace{I \otimes I \otimes \cdots \otimes I})$, which is the following:

$$
\rho_{2 d}(n):=\min _{r \in \mathbb{Z}_{+}}\left\{\exists b^{1}, b^{2}, \ldots, b^{r} \in \mathbb{R}^{n}, \lambda \in \mathbb{R}^{r}, \text { such that }\left(x^{\mathrm{T}} x\right)^{d}=\sum_{i=1}^{r} \lambda_{i}\left(\left(b^{i}\right)^{\mathrm{T}} x\right)^{2 d} \forall x \in \mathbb{R}^{n}\right\},
$$

or in the language of tensors, the smallest $r$ such that

$$
\operatorname{sym}(\underbrace{I \otimes I \otimes \cdots \otimes I}_{d})=\sum_{i=1}^{r} \lambda_{i} \underbrace{b^{i} \otimes b^{i} \otimes \cdots \otimes b^{i}}_{d}
$$

The difference between $\tau_{2 d}(n)$ and $\rho_{2 d}(n)$ lies in the fact that the latter one allows negative rank-one tensors. Therefore we have $\tau_{2 d}(n) \geq \rho_{2 d}(n)$. Computing the exact values for $\tau_{2 d}(n)$ and $\rho_{2 d}(n)$ is not easy for general $n$ and $d$, and the only clear case is for $d=1$ whereas $\tau_{2}(n)=\rho_{2}(n)=n$. In this section, we focus on the case $d=2$, i.e., $\tau_{4}(n)$ and $\rho_{4}(n)$. In fact, the lower bound for $\tau_{2 d}(n)$ was already studied by Reznick [17]. Below we first summarize the result of Reznick [17].

Theorem 5.1 (Theorem 8.15 of Reznick [17]). For any given positive integers $d$ and $n$, the number of $d$-th powered linear terms in Hilberts identity (5) is at least $\binom{n+d-1}{n-1}$, i.e., $\tau_{2 d}(n) \geq\binom{ n+d-1}{n-1}$.

Furthermore when $d=2$, the exact values $\tau_{2 d}(n)$ for some specific $n$ 's are known in the literature.
Proposition 5.1 (Proposition 9.26 of Reznick [17]).

$$
\tau_{4}(n)=\binom{n+2-1}{n-1}+1=\frac{n(n+1)}{2}+1 \quad \text { when } n=4,5,6 .
$$

We remark that when $d=2, n(n+1) / 2$ is also a lower bound for the number of rank-one terms to represent $\operatorname{sym}(A \otimes A)$ with $A \succ 0$. Besides, if $\xi_{1}, \xi_{2}, \ldots, \xi_{n}$ are symmetric Bernoulli random variables, and they are 4-wise uncorrelated, then Theorem 5.1 also indicates that $n(n+1) / 2$ is a lower bound for the size of sample space generated by $\left\{\xi_{1}, \xi_{2}, \ldots, \xi_{n}\right\}$. In fact, $n(n+1) / 2$ is also a lower bound for the rank of $\operatorname{sym}(I \otimes I)$, as the following theorem stipulates.

Theorem 5.2. For any positive integer $n$, it holds that $n(n+1) / 2 \leq \rho_{4}(n) \leq n^{2}$.
Proof. Denote the shortest representation to be

$$
\left(\sum_{j=1}^{n} x_{j}^{2}\right)^{2}=\sum_{i=1}^{m}\left(\sum_{j=1}^{n} a_{i j} x_{j}\right)^{4}-\sum_{i=1}^{l}\left(\sum_{j=1}^{n} b_{i j} x_{j}\right)^{4},
$$

where $m+l=\rho_{4}(n)$. By comparing the coefficient of each monomial, we have

$$
\begin{gather*}
\sum_{i=1}^{m} a_{i j}^{4}-\sum_{i=1}^{l} b_{i j}^{4}=1 \quad \forall 1 \leq j \leq n \\
\sum_{i=1}^{m} a_{i j_{1}}^{2} a_{i j_{2}}^{2}-\sum_{i=1}^{l} b_{i j_{1}}^{2} b_{i j_{2}}^{2}=\frac{1}{3} \quad \forall 1 \leq j_{1} \neq j_{2} \leq n \\
\sum_{i=1}^{m} a_{i j_{1}}^{3} a_{i j_{2}}-\sum_{i=1}^{l} b_{i j_{1}}^{3} b_{i j_{2}}=0 \quad \forall 1 \leq j_{1} \neq j_{2} \leq n  \tag{19}\\
\sum_{i=1}^{m} a_{i j_{1}}^{2} a_{i j_{2}} a_{i j_{3}}-\sum_{i=1}^{l} b_{i j_{1}}^{2} b_{i j_{2}} b_{i j_{3}}=0 \quad \forall 1 \leq j_{1}, j_{2}, j_{3} \leq n \text { with } j_{k} \neq j_{t} \text { if } k \neq t \\
\sum_{i=1}^{m} a_{i j_{1}} a_{i j_{2}} a_{i j_{3}} a_{i j_{4}}-\sum_{i=1}^{l} b_{i j_{1}} b_{i j_{2}} b_{i j_{3}} b_{i j_{4}}=0 \quad \forall 1 \leq j_{1}, j_{2}, j_{3}, j_{4} \leq n \text { with } j_{k} \neq j_{t} \text { if } k \neq t
\end{gather*}
$$

Construct matrices $A \in \mathbb{R}^{m \times n}, C \in \mathbb{R}^{l \times n}, B \in \mathbb{R}^{m \times(n(n-1)) / 2}$, and $D \in \mathbb{R}^{l \times(n(n-1)) / 2}$, where

$$
\begin{aligned}
& A=\left[\begin{array}{cccc}
a_{11}^{2} & a_{12}^{2} & \ldots & a_{1 n}^{2} \\
a_{21}^{2} & a_{22}^{2} & \ldots & a_{2 n}^{2} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1}^{2} & a_{m 2}^{2} & \ldots & a_{m n}^{2}
\end{array}\right], \quad C=\left[\begin{array}{cccc}
b_{11}^{2} & b_{12}^{2} & \ldots & b_{1 n}^{2} \\
b_{21}^{2} & b_{22}^{2} & \ldots & b_{2 n}^{2} \\
\vdots & \vdots & \ddots & \vdots \\
b_{l 1}^{2} & b_{l 2}^{2} & \ldots & b_{l n}^{2}
\end{array}\right], \\
& B=\left[\begin{array}{cccccccccc}
a_{11} a_{12} & a_{11} a_{13} & \ldots & a_{11} a_{1 n} & a_{12} a_{13} & a_{12} a_{14} & \ldots & a_{12} a_{1 n} & \ldots & a_{1, n-1} a_{1 n} \\
a_{21} a_{22} & a_{21} a_{23} & \ldots & a_{21} a_{2 n} & a_{22} a_{23} & a_{22} a_{24} & \ldots & a_{22} a_{2 n} & \ldots & a_{2, n-1} a_{2 n} \\
\vdots & \vdots & \ddots & \vdots & \vdots & \vdots & \ddots & \vdots & \ddots & \vdots \\
a_{m 1} a_{m 2} & a_{m 1} a_{m 3} & \ldots & a_{m 1} a_{m n} & a_{m 2} a_{m 3} & a_{m 2} a_{m 4} & \ldots & a_{m 2} a_{m n} & \ldots & a_{m, n-1} a_{m n}
\end{array}\right],
\end{aligned}
$$

and

$$
D=\left[\begin{array}{cccccccccc}
b_{11} b_{12} & b_{11} b_{13} & \ldots & b_{11} b_{1 n} & b_{12} b_{13} & b_{12} b_{14} & \ldots & b_{12} b_{1 n} & \ldots & b_{1, n-1} b_{1 n} \\
b_{21} b_{22} & b_{21} b_{23} & \ldots & b_{21} b_{2 n} & b_{22} b_{23} & b_{22} b_{24} & \ldots & b_{22} b_{2 n} & \ldots & b_{2, n-1} b_{2 n} \\
\vdots & \vdots & \ddots & \vdots & \vdots & \vdots & \ddots & \vdots & \ddots & \vdots \\
b_{l 1} b_{l 2} & b_{l 1} b_{l 3} & \ldots & b_{l 1} b_{l n} & b_{l 2} b_{l 3} & b_{l 2} b_{l 4} & \ldots & b_{l 2} b_{l n} & \ldots & b_{l, n-1} b_{l n}
\end{array}\right] .
$$

By (19), it is straightforward to verify that

$$
[A, B]^{\mathrm{T}}[A, B]-[C, D]^{\mathrm{T}}[C, D]=\left[\begin{array}{cc}
A^{\mathrm{T}} A-C^{\mathrm{T}} C & A^{\mathrm{T}} B-C^{\mathrm{T}} D \\
B^{\mathrm{T}} A-D^{\mathrm{T}} C & B^{\mathrm{T}} B-D^{\mathrm{T}} D
\end{array}\right]=\left[\begin{array}{cc}
\frac{1}{3} E+\frac{2}{3} I & O \\
O & \frac{1}{3} I
\end{array}\right] \succ 0
$$

Thus $[A, B]^{\mathrm{T}}[A, B]$ is also positive definite, hence full rank. Finally,

$$
\rho_{4}(n) \geq m \geq \operatorname{rank}([A, B]) \geq \operatorname{rank}\left([A, B]^{\mathrm{T}}[A, B]\right)=n(n+1) / 2
$$

The upper bound follows from the following identity (Reznick [17, (10.35)]):

$$
\left(\sum_{j=1}^{n} x_{j}^{2}\right)^{2}=\frac{1}{6} \sum_{j<k}\left(x_{j}+x_{k}\right)^{4}+\frac{1}{6} \sum_{j<k}\left(x_{j}-x_{k}\right)^{4}+\frac{4-n}{3} \sum_{j=1}^{n} x_{j}^{4} .
$$

When $n \geq 5$, the coefficient $(4-n) / 3$ is negative, and so it is not a valid representation of Hilbert's identity, but it is still a rank-one decomposition for $\left(\sum_{j=1}^{n} x_{j}^{2}\right)^{2}$. Since there are no more than $n^{2}$ rank-one terms in this expression, it yields an upper bound of $n^{2}$ for $\rho_{4}(n)$.

Remark that as $\rho_{4}(n) \leq \tau_{4}(n)$, Theorem 5.2 immediately implies Theorem 5.1 when $d=2$. The following examples show that $n(n+1) / 2$ is the exact value for $\rho_{4}(n)$ as well as $\tau_{4}(n)$ when $n \leq 3$ (note that the case $n=1$ is trivial).

Example 5.1. $\quad\left(x_{1}^{2}+x_{2}^{2}\right)^{2}=\frac{1}{2}\left(x_{1}+(1 / \sqrt{3}) x_{2}\right)^{4}+\frac{1}{2}\left(x_{1}-(1 / \sqrt{3}) x_{2}\right)^{4}+\frac{8}{9} x_{2}^{4}$.
Example 5.2. $\quad\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}\right)^{2}=\left(1 /\left(2\left(a^{4}+1\right)\right)\right) \sum_{\beta= \pm 1}\left(\left(x_{1}+\beta a x_{2}\right)^{4}+\left(x_{2}+\beta a x_{3}\right)^{4}+\left(x_{3}+\beta a x_{1}\right)^{4}\right)$, where $a^{2}=(3 \pm \sqrt{5}) / 2$.

We remark that the above tight representations are not unique. One may find other representations, e.g., Reznick [17, (8.29) and (8.30)], which include Examples 5.1 and 5.2 as special cases. Moreover, in light of Proposition 5.1, Liouville's identity (6), which involving 12 terms, is not tight for both $\rho_{4}(4)$ and $\tau_{4}(4)$. The following tight example for $\tau_{4}(4)$ only includes 11 terms.

Example 5.3 (Reznick [17, (9.27)(i)]).

$$
\begin{aligned}
\left(x_{1}^{2}+x_{2}^{2}+x_{3}^{2}+x_{4}^{2}\right)^{2}= & \frac{1}{32}\left(x_{1}+x_{2}+x_{3}+x_{4}\right)^{4}+\frac{1}{192} \sum_{i=1}^{4}\left(3 x_{i}-\sum_{1 \leq j \leq 4,, j \neq i} x_{j}\right)^{4} \\
& +\frac{1}{192} \sum_{1 \leq i<j \leq 4}\left((1+\sqrt{2})\left(x_{i}+x_{j}\right)-(1-\sqrt{2}) \sum_{1 \leq k \leq 4, k \neq i, j} x_{k}\right)^{4} .
\end{aligned}
$$

This example along with Theorem 5.2 implies that $10 \leq \rho_{4}(4) \leq 11$. It remains an open problem to locate the exact value of $\rho_{4}(4)$. In general, finding the exact values (or a tighter upper bound) of $\tau_{4}(n)$ and $\rho_{4}(n)$, as well as finding a succinct algorithm to construct a shorter (less than $2 n^{4}+n$ ) representation of the fourth moments tensor $\operatorname{sym}(I \otimes I)$, are interesting future research questions.
6. Matrix $q \mapsto p$ norm problem. In this section, we shall illustrate the power of polynomial-size representation of moments tensor by a specific example. In particular, we consider the problem of computing the so-called $q \mapsto p(1 \leq p, q \leq \infty)$ norm of a matrix $A$, defined as follows:

$$
\|A\|_{q \mapsto p}:=\max _{\|x\|_{q}=1}\|A x\|_{p}
$$

This problem can be viewed as a natural extension of several useful problems. For instance, the case $p=q=2$ corresponds to the largest singular value of $A$. The case $(p, q)=(1, \infty)$ corresponds to the bilinear form optimization problem in binary variables, which is related to the so-called matrix cut norm and Grothendieck's constant; see Alon and Naor [1]. In case $p=q$, the problem becomes the matrix $p$-norm problem, which has applications in scientific computing; cf. Higham [9].

In terms of the computational complexity, three easy cases are well known: (1) $q=1$ and $p \geq 1$ is a rational number; (2) $p=\infty$ and $q \geq 1$ is a rational number; (3) $p=q=2$. Steinberg [19] showed that computing $\|A\|_{q \mapsto p}$ is NP-hard for general $1 \leq p<q \leq \infty$, and she further conjectured that the above mentioned three cases are the only exceptional easy cases where the matrix $q \mapsto p$ norm can be computed in polynomial time. Hendrickx and Olshevsky [8] made some progress along this line by figuring out the complexity status of the "diagonal" case of $p=q$. Moreover, very recently Bhaskara and Vijayaraghavan [5] proved that this problem is NP-hard to approximate to any constant factor when $2<p \leq q$. However, the problem of determining the complexity status for the case $p>q$ still remains open. Here we shall show that the problem $\|A\|_{q \mapsto p}$ is NP-hard when $p=4$ and $q=2$. To this end, let us first present the following lemma.

Lemma 6.1. Given positive integers $n, i, j$ with $1 \leq i<j \leq n$, we can find $t\left(\leq 2 n^{4}+n+2\right)$ vectors $a^{1}, a^{2}$, $\ldots, a^{t}$ in polynomial time, such that

$$
2 x_{i}^{2} x_{j}^{2}+\left(x^{\mathrm{T}} x\right)^{2}=\sum_{k=1}^{t}\left(\left(a^{k}\right)^{\mathrm{T}} x\right)^{4}
$$

Proof. Recall in Theorem 4.1, we can find $\tau\left(\leq 2 n^{4}\right)$ vectors $a^{1}, a^{2}, \ldots, a^{\tau} \in \mathbb{R}^{n}$ in polynomial time, such that

$$
\begin{equation*}
\frac{2}{3} \sum_{l=1}^{n} x_{l}^{4}+\sum_{l=1}^{\tau}\left(\left(a^{l}\right)^{\mathrm{T}} x\right)^{4}=\left(x^{\mathrm{T}} x\right)^{2} \tag{20}
\end{equation*}
$$

On the other hand, one verifies straightforwardly that for $1 \leq i \neq j \leq n$ we have

$$
\begin{equation*}
\frac{1}{2}\left(\left(x_{i}+x_{j}\right)^{4}+\left(x_{i}-x_{j}\right)^{4}\right)+x_{i}^{4}+x_{j}^{4}+2 \sum_{1 \leq l \leq n, l \neq i, j} x_{l}^{4}=6 x_{i}^{2} x_{j}^{2}+2 \sum_{l=1}^{n} x_{l}^{4} \tag{21}
\end{equation*}
$$

Dividing by 3 on both sides of (21) and then summing up with $\sum_{l=1}^{\tau}\left(\left(a^{l}\right)^{\mathrm{T}} x\right)^{4}$ yields

$$
\begin{aligned}
& \sum_{l=1}^{\tau}\left(\left(a^{l}\right)^{\mathrm{T}} x\right)^{4}+\frac{1}{3}\left(\frac{1}{2}\left(\left(x_{i}+x_{j}\right)^{4}+\left(x_{i}-x_{j}\right)^{4}\right)+x_{i}^{4}+x_{j}^{4}+2 \sum_{1 \leq l \leq n, l \neq i, j} x_{l}^{4}\right) \\
& \quad=\sum_{l=1}^{\tau}\left(\left(a^{l}\right)^{\mathrm{T}} x\right)^{4}+2 x_{i}^{2} x_{j}^{2}+\frac{2}{3} \sum_{l=1}^{n} x_{l}^{4} \\
& \quad=2 x_{i}^{2} x_{j}^{2}+\left(x^{\mathrm{T}} x\right)^{2}
\end{aligned}
$$

where the last equality is due to (20).
Now we are in a position to prove the main theorem of this section.
Theorem 6.1. Computing $\|A\|_{2 \mapsto 4}=\max _{\|x\|_{2}=1}\|A x\|_{4}$ is NP-hard.
Proof. The reduction is made from computing the maximum (vertex) independence set of a graph. In particular, for a given graph $G=(V, E)$, Nesterov [15] showed that the following problem can be reduced from the maximum independence number problem:

$$
\begin{aligned}
\max & 2 \sum_{(i, j) \in E, i<j} x_{i}^{2} x_{j}^{2} \\
\text { s.t. } & \|x\|_{2}=1, \quad x \in \mathbb{R}^{n},
\end{aligned}
$$

hence is NP-hard. Moreover, the above is obviously equivalent to

$$
\begin{align*}
\max & 2 \sum_{(i, j) \in E, i<j} x_{i}^{2} x_{j}^{2}+|E| \cdot\|x\|_{2}^{4}=\sum_{(i, j) \in E, i<j}\left(2 x_{i}^{2} x_{j}^{2}+\left(x^{\mathrm{T}} x\right)^{2}\right)  \tag{P}\\
\text { s.t. } & \|x\|_{2}=1, \quad x \in \mathbb{R}^{n} .
\end{align*}
$$

By Lemma 6.1, the objective in (P) can be expressed by no more than $|E| \cdot\left(2 n^{4}+n+2\right)$ number of fourth powered linear terms, making (P) be an instance of $\|A\|_{2 \mapsto 4}$ (polynomial size). The polynomial reduction is thus complete.

Suppose that $p^{\prime}$ and $q^{\prime}$ are the conjugates of $p$ and $q$, respectively, i.e., $1 / p+1 / p^{\prime}=1$ and $1 / q+1 / q^{\prime}=1$. By using the fact that $\|x\|_{p}=\max _{\|y\|_{p^{\prime}}=1} y^{\mathrm{T}} x$, one can prove that $\|A\|_{q \mapsto p}=\left\|A^{\mathrm{T}}\right\|_{p^{\prime} \mapsto q^{\prime}}$. Therefore, Theorem 6.1 implies that computing $\|A\|_{4 / 3 \mapsto 2}$ is also NP-hard. We remark that Theorem 6.1 was independently proved by Barak et al. [3] using a similar argument, after the initial version of this paper was submitted.

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[^0]:    ${ }^{1}$ During the review process of this paper, Barak et al. [3] independently proved that it is NP-hardness to compute the matrix $2 \mapsto 4$ norm.

