# On the numerical condition of a generalized Hankel eigenvalue problem 

by

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#### Abstract

The generalized eigenvalue problem $\widetilde{H} y=\lambda H y$ with $H$ a Hankel matrix and $\widetilde{H}$ the corresponding shifted Hankel matrix occurs in number of applications such as the reconstruction of the shape of a polygon from its moments, the determination of abscissa of quadrature formulas, of poles of Padé approximants, or of the unknown powers of a sparse black box polynomial in computer algebra. In many of these applications, the entries of the Hankel matrix are only known up to a certain precision. We study the sensitivity of the nonlinear application mapping the vector of Hankel entries to its generalized eigenvalues.

A basic tool in this study is a result on the condition number of Vandermonde matrices with not necessarily real abscissas which are possibly row-scaled.


Key words: Shape detection via moments, ill-disposed eigenvalues, condition number, Hankel matrix, Vandermonde matrix.
AMS Classification (2000): 15A18, 65F35, 15A12, 30E10.

## 1 Introduction

Given complex numbers $h_{0}, h_{1}, \ldots, h_{2 n-1}$, which we refer to as moments, we are interested in the generalized eigenvalue problem (abbreviated as GEP)

$$
\begin{equation*}
\widetilde{H}_{n} y^{R}=\lambda H_{n} y^{R}, \quad y^{L} \widetilde{H}_{n}=\lambda y^{L} H_{n} \tag{1}
\end{equation*}
$$

for the Hankel matrices

$$
H_{n}:=\left[\begin{array}{cccc}
h_{0} & h_{1} & \cdots & h_{n-1} \\
h_{1} & h_{2} & \cdots & h_{n} \\
\vdots & \vdots & & \vdots \\
h_{n-1} & h_{n} & \cdots & h_{2 n-2}
\end{array}\right], \quad \widetilde{H}_{n}:=\left[\begin{array}{cccc}
h_{1} & h_{2} & \cdots & h_{n} \\
h_{2} & h_{3} & \cdots & h_{n+1} \\
\vdots & \vdots & & \vdots \\
h_{n} & h_{n+1} & \cdots & h_{2 n-1}
\end{array}\right] .
$$

Such eigenvalue problems appear in applications such as shape reconstruction of a polygon from its moments [GMV99], determination of abscissa of quadrature formulas [Gau82, Gau86] and the determination of the unknown powers of a multivariate sparse black box polynomial in computer algebra [GLL04a, GLL06]. The shape from moments problem is applicable to a wide variety of inverse problems of uniform density regions related to general elliptical equations appearing for example in statistics and probability, geophysics, and computed tomography (see for instance the references in [GMV99]). In the case of sparse multivariate interpolation the applications

[^0]include approximate multivariate factorization and decompositions of approximately specified polynomial systems (see the references in [GLL06]).

In many applications, the moments $h_{k}$ are obtained by measurements and hence are affected by some error. In this paper we are interested in the sensitivity of the generalized eigenvalues with respect to such errors in the moments. In other words, we are interested in the sensitivity of the application

$$
\begin{align*}
G_{n}: & \mathbb{C}^{2 n} \mapsto \mathbb{C}^{n} \\
& \left(h_{0}, \ldots, h_{2 n-1}\right) \mapsto\left(\lambda_{1}, \ldots, \lambda_{n}\right) \tag{2}
\end{align*}
$$

which maps the moments to the $n$ generalized eigenvalues defined by (1). Up to first order, this sensitivity is measured by the norm of the Jacobian of this non-linear map. In the case $\lambda_{j} \in \mathbb{R}$ and $c_{j}>0$, which is related to orthogonality on the real line, such a study was done previously by Gautschi [Gau82, Gau86] followed by other authors [Fis96, BeBo98]. In our applications, however, we are more interested in the case where $\lambda_{j}$ and $c_{j}$ are complex.

In the applications mentioned earlier the input is a set of given moments $h_{k}$ satisfying

$$
\begin{equation*}
h_{k}=\sum_{j=1}^{n} c_{j} \lambda_{j}^{k}, \quad k=0, \ldots, 2 n-1, \tag{3}
\end{equation*}
$$

with $c_{j} \in \mathbb{C} \backslash\{0\}$ and distinct $\lambda_{j} \in \mathbb{C}$. In such applications, the quantities $c_{j}$ and $\lambda_{j}$ are unknown, and need to be computed from the moments. Notice that the $\lambda_{j}$ in formula (3) are just the generalized eigenvalues of problem (1). Conversely, one may in fact show that if (1) has distinct generalized eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$ then there exist suitable $c_{j}$ such that equation (3) holds true. Thus, the generalized eigenvalue problem (1) is a convenient way of writing the problem of finding the $\lambda_{j}$ from the moments $h_{k}$. Additional details for the above points are given later in §2.1.

This paper has two main contributions. The first is the (rather surprising) result that the sensitivity of (1) with respect to both structured and unstructured perturbations can be measured essentially by the same quantity. This gives a theoretical justification for a numerical observation in [GMV99]. Our second contribution is to give lower bounds for the (relative) sensitivity for the most sensitive eigenvalue of (1). These lower bounds are given in terms of the condition number of the underlying Hankel matrix $H_{n}$. For a subclass of such problems (which are relevant to the applications mentioned earlier) there are lower bounds for the most sensitive eigenvalue that is given in terms of an associated Vandermonde matrix formed with the help of the (possibly complex) $\lambda_{j}$ multiplied by some diagonal factor. In order to obtain more precise information about the sensitivity in terms of these $\lambda_{j}$, we finally analyze more closely the condition number of such scaled Vandermonde matrices. Our analysis shows that this sensitivity potentially grows exponentially with the size of the matrix.

The remainder of the paper is organized as follows. In $\S 2$ we recall the notion of ill-disposed generalized eigenvalues from [GMV99] and compare the sensitivity of (1) with respect to both structured and unstructured perturbations. In $\S 3$ we study the sensitivity for the most sensitive eigenvalue of (1) and its relation to the condition number of the underlying Hankel matrix. In this section we also give lower bounds for the most sensitive eigenvalue for a subclass of problems relevant for our applications, and study the dependency of this lower bound on the distribution of the $\lambda_{j}$ in the complex plane. In $\S 4$ we illustrate the impact of our results by discussing some simple examples coming from the applications of shape detection and the reconstruction of a sparse black box polynomial. The proof of the main Theorem of $\S 3$ is given in $\S 5$. Finally, a conclusion and discussion for further work is given in $\S 6$.

Notations: We let $\|\cdot\|$ be the Euclidean vector norm and the spectral matrix norm and $\|S\|_{F}=\sqrt{\operatorname{trace}\left(S^{*} S\right)}$ the subordinate Frobenius matrix norm, where $S^{*}$ denotes the adjoint of the matrix $S$. The $j$-th canonical vector is denoted by $e_{j}$. For a continuous function $f: \mathbb{C} \mapsto \mathbb{C}$, and some compact set $E \subset \mathbb{C}$, we consider the norms

$$
\|f\|_{L_{2}(\partial \mathbb{D})}:=\left[\frac{1}{2 \pi} \int_{0}^{2 \pi}\left|f\left(e^{i t}\right)\right|^{2} d t\right]^{1 / 2}, \quad\|f\|_{L_{\infty}(E)}:=\max _{z \in E}|f(z)|,
$$

where $\mathbb{D}$ denotes the closed unit disk, and $\partial \mathbb{D}$ its boundary. Given a polynomial $p(z)=a_{0}+$ $a_{1} z+\ldots+a_{n} z^{n}$, we consider its vector of coefficients $\vec{p}:=\left(a_{0}, a_{1}, \ldots, a_{n}\right)^{t}$, the length of this vector being clear from the context. One easily verifies that, for any polynomial $p$ of degree at most $n-1$, the following bounds hold

$$
\begin{equation*}
\operatorname{deg} p<n: \quad\|\vec{p}\|=\|p\|_{L_{2}(\partial \mathbb{D})}, \quad \frac{1}{\sqrt{n}}\|p\|_{L_{\infty}(\partial \mathbb{D})} \leq\|p\|_{L_{2}(\partial \mathbb{D})} \leq\|p\|_{L_{\infty}(\partial \mathbb{D})} \tag{4}
\end{equation*}
$$

## 2 Unstructured and structured perturbations of the GEP

In this section we give some background for the generalized eigenvalue problem (1). In particular, we show that in the generic case it can always be formulated in terms of problem (3) and give a simple well-known formula for the generalized eigenvectors. We also recall the perturbation analysis of this problem for the case of unstructured perturbations from [GMV99] and then give the main result of this section, a perturbation analysis which takes into account the special Hankel structure of our input matrices.

### 2.1 Preliminaries

Let us first show that the generalized eigenvalue problem (1) is equivalent to finding $\lambda_{j}$ in equation (3) from the moments $h_{k}$ in the generic case. Thus, suppose problem (1) has moments defined by (3). Then one has a simple formula for the corresponding generalized eigenvalues and eigenvectors. Indeed in this case one has the factorizations of the two Hankel matrices

$$
\begin{equation*}
H_{n}=V_{n}^{t} \operatorname{diag}\left(c_{1}, \ldots, c_{n}\right) V_{n}=W_{n}^{t} W_{n}, \quad \widetilde{H}_{n}=V_{n}^{t} \operatorname{diag}\left(\lambda_{1} c_{1}, \ldots, \lambda_{n} c_{n}\right) V_{n}, \tag{5}
\end{equation*}
$$

where ${ }^{t}$ denotes taking of a transpose (without taking conjugates), and

$$
V_{n}=\left[\begin{array}{cccc}
1 & \lambda_{1} & \cdots & \lambda_{1}^{n-1}  \tag{6}\\
1 & \lambda_{2} & \cdots & \lambda_{2}^{n-1} \\
\vdots & \vdots & & \vdots \\
1 & \lambda_{n} & \cdots & \lambda_{n}^{n-1}
\end{array}\right], \quad W_{n}:=\operatorname{diag}\left(\sqrt{c_{1}}, \ldots, \sqrt{c_{n}}\right) V_{n},
$$

a Vandermonde matrix and a row-scaled Vandermonde matrix, respectively. It follows from (5) that the generalized eigenvalues of (1) are given by the abscissa $\lambda_{j}$ of (3), with corresponding right and left eigenvectors given by $y^{R}=\left(y^{L}\right)^{t}=V_{n}^{-1} e_{j}$.

Conversely suppose that (1) has distinct generalized eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$. Then it is wellknown [BaGM95] that the polynomial $q(z):=\operatorname{det}\left(\widetilde{H}_{n}-z H_{n}\right)$ is the denominator of the $n$-th Padé approximant at infinity of the series $h(z)=h_{0} z^{-1}+h_{1} z^{-2}+\cdots$. Thus the generalized eigenvalues of (1) are just the poles of the $n$-th Padé approximant. In this case we have a partial fraction decomposition for the $n$-th Padé approximant given by

$$
\sum_{j=0}^{2 n-1} \frac{h_{j}}{z^{j+1}}=\sum_{j=1}^{n} \frac{c_{j}}{z-\lambda_{j}}+\mathcal{O}\left(z^{-2 n-1}\right)_{z \rightarrow \infty},
$$

with $c_{j} \in \mathbb{C} \backslash\{0\}$. Equating coefficients leads to the representation (3).

### 2.2 Analysis of Errors

In [GMV99, Section 3.3], Golub, Milanfar and Varah considered perturbations of the generalized eigenvalue problem (1), namely

$$
\begin{equation*}
\left[\widetilde{H}_{n}+\epsilon \widetilde{E}_{n}\right] y^{R}(\epsilon)=\lambda(\epsilon)\left[H_{n}+\epsilon E_{n}\right] y^{R}(\epsilon), \quad y^{L}(\epsilon)\left[\widetilde{H}_{n}+\epsilon \widetilde{E}_{n}\right]=\lambda(\epsilon) y^{L}(\epsilon)\left[H_{n}+\epsilon E_{n}\right], \tag{7}
\end{equation*}
$$

where $\epsilon>0$ is small, and the perturbation matrices $\widetilde{E}_{n}, E_{n}$ normalized such that $\left\|\tilde{E}_{n}\right\| \leq 1$, $\left\|E_{n}\right\| \leq 1$. In a small neigborhood around a simple eigenvalue $\lambda(0)=\lambda_{j}$ (and thus $y^{R}(0)=$ $\left.\left(y^{L}(0)\right)^{t}=V_{n}^{-1} e_{j}\right)$, the function $\epsilon \mapsto \lambda(\epsilon)=\lambda_{j}(\epsilon)$ is differentiable, with its derivative at zero being given by

$$
\frac{d \lambda_{j}}{d \epsilon}(0)=\frac{y^{L}(0)\left[\tilde{E}_{n}-\lambda_{j} E_{n}\right] y^{R}(0)}{y^{L}(0) H_{n} y^{R}(0)} .
$$

Plugging in explicitly the eigenvectors and using the first formula of (5) leads to the expression

$$
\begin{equation*}
\frac{d \lambda_{j}}{d \epsilon}(0)=\frac{\left(V_{n}^{-1} e_{j}\right)^{t}\left[\widetilde{E}_{n}-\lambda_{j} E_{n}\right] V_{n}^{-1} e_{j}}{\left(V_{n}^{-1} e_{j}\right)^{t} H_{n} V_{n}^{-1} e_{j}}=\frac{\left(V_{n}^{-1} e_{j}\right)^{t}\left[\widetilde{E}_{n}-\lambda_{j} E_{n}\right] V_{n}^{-1} e_{j}}{c_{j}} . \tag{8}
\end{equation*}
$$

Following [GMV99, Section 3.3],we say that a generalized eigenvalue $\lambda_{j}$ is ill-disposed if the right-hand side of (8) multiplied with $\left\|H_{n}\right\|+\left\|\widetilde{H}_{n}\right\|$ is "large".

Using the link (6) between $V_{n}$ and $W_{n}$, and exploiting the given information on the norms of $E_{n}$ and $\widetilde{E}_{n}$, we obtain from (8) the simpler upper bound

$$
\begin{equation*}
\left|\frac{d \lambda_{j}}{d \epsilon}(0)\right| \leq\left(\left|\lambda_{j}\right|+1\right)\left\|W_{n}^{-1} e_{j}\right\|^{2} \tag{9}
\end{equation*}
$$

This upper bound appears to be quite rough. In addition, the bound is based on unstructured perturbations of Hankel matrices which a priori should lead to a severe over estimation of the actual sensitivity since in applications the perturbations will be structured. For instance, if we are interested in the sensitivity of $\lambda_{j}$ with respect to small perturbations of the second moment $h_{2}$, we should consider as perturbation matrices $E_{n}, \widetilde{E}_{n}$ those matrices obtained from $H_{n}, \widetilde{H}_{n}$ by replacing $h_{2}$ by 1 , and all other moments by 0 .

However, it follows from Theorem 2.1 below that the sensitivity with respect to structured perturbations essentially gives the same result as (9). This fact was already observed numerically in [GMV99]. To make a clear statement about first order sensitivity of the generalized eigenvalues $\lambda_{j}$ with respect to perturbations on the moments $h_{k}$, we have to compute the Jacobian of the non-linear map $G_{n}$ defined in (2), and to estimate the norms of its rows.

Theorem 2.1. Suppose that (1) has distinct possibly complex eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$, and let $c_{j}$, and $W_{n}$, as in (3), and (6), respectively. Then for all $j=1, \ldots, n$ we have

$$
\left\|\left(\frac{\partial \lambda_{j}}{\partial h_{k}}\right)_{k=0, \ldots, 2 n-1}\right\|=\eta_{j, n}\left(\left|\lambda_{j}\right|+1\right)\left\|W_{n}^{-1} e_{j}\right\|^{2}
$$

where $1 /(2 n) \leq \eta_{j, n} \leq \sqrt{2 n}$.

Proof. We start with the general observation that polynomial language is very useful for express the inverse of a Vandermonde matrix $V_{n}$ with abscissa $\lambda_{1}, \ldots, \lambda_{n}$. Indeed, with the corresponding

Lagrange polynomials

$$
\begin{equation*}
\ell_{j}(z)=\prod_{k=1, \ldots, n, k \neq j} \frac{z-\lambda_{k}}{\lambda_{j}-\lambda_{k}}=\frac{\omega(z)}{\left(z-\lambda_{j}\right) \omega^{\prime}\left(\lambda_{j}\right)}, \quad \omega(z)=\prod_{k=1}^{n}\left(z-\lambda_{k}\right) \tag{10}
\end{equation*}
$$

we get the relation $\ell_{j}\left(\lambda_{k}\right)=\delta_{j, k}$ or, in other words,

$$
\begin{equation*}
j=1, \ldots, n: \quad V_{n}^{-1} e_{j}=\vec{\ell}_{j}, \quad\left\|V_{n}^{-1} e_{j}\right\|=\left\|\ell_{j}\right\|_{L_{2}(\partial \mathbb{D})} \tag{11}
\end{equation*}
$$

Consider a fixed $j \in\{1, . ., n\}$ and a fixed $k \in\{0,1, \ldots, 2 n-1\}$. In order to measure the first order sensitivity of the generalized eigenvalue $\lambda_{j}$ with respect to perturbations in the moment $h_{k}$, we proceed as in (8) and obtain the formula

$$
\frac{\partial \lambda_{j}}{\partial h_{k}}=\frac{\left(V_{n}^{-1} e_{j}\right)^{t}\left[\tilde{E}_{n}-\lambda_{j} E_{n}\right] V_{n}^{-1} e_{j}}{c_{j}}
$$

However, now the perturbation matrices $E_{n}$ and $\tilde{E}_{n}$ have a special form, namely the matrix $E_{n}$ (respectively $\tilde{E}_{n}$ ) is the Hankel matrix obtained from the zero matrix by replacing the $(k+1)$-th anti-diagonal by ones (respectively the $k$-th anti-diagonal). Again it is useful to express this formula in polynomial language. With the corresponding Lagrange polynomials $\ell_{j}(z)=\sum_{\kappa} \ell_{j, \kappa} z^{\kappa}$ as in (10) and the two polynomials

$$
p(z)=\ell_{j}(z)^{2}=\sum_{\kappa} p_{\kappa} z^{\kappa}, \quad q(z)=\left(z-\lambda_{j}\right) \ell_{j}(z)^{2}=\sum_{\kappa} q_{\kappa} z^{\kappa}
$$

we get from (11) that

$$
\begin{aligned}
& \left(V_{n}^{-1} e_{j}\right)^{t} E_{n} V_{n}^{-1} e_{j}=\sum_{m=0}^{k} \ell_{m, j} \ell_{k-m, j}=p_{k} \\
& \left(V_{n}^{-1} e_{j}\right)^{t} \tilde{E}_{n} V_{n}^{-1} e_{j}=\sum_{m=0}^{k} \ell_{m, j} \ell_{k-1-m, j}=p_{k-1}
\end{aligned}
$$

and thus

$$
\frac{\partial \lambda_{j}}{\partial h_{k}}=\frac{1}{c_{j}}\left[p_{k-1}-\lambda_{j} p_{k}\right]=\frac{q_{k}}{c_{j}}
$$

Using different techniques, a similar formula has been found by Gautschi (see [Gau82, Section 3.2] and [Gau86]) for the case $\lambda_{j} \in \mathbb{R}$ and $c_{j}>0$.

Since $W_{n}^{-1} e_{j}=V_{n}^{-1} e_{j} / \sqrt{c_{j}}=\vec{\ell}_{j} / \sqrt{c_{j}}$ by (6) and (11), in order to establish Theorem 2.1 it only remains to show that

$$
\begin{equation*}
\frac{\left|\lambda_{j}\right|+1}{2 n}\left\|\vec{\ell}_{j}\right\|^{2} \leq\|\vec{q}\| \leq \sqrt{2 n}\left(\left|\lambda_{j}\right|+1\right)\left\|\vec{\ell}_{j}\right\|^{2} \tag{12}
\end{equation*}
$$

Notice that $\vec{q}=B \vec{p}$, with $B$ being a matrix of size $(2 n) \times(2 n-1)$ obtained from the zero matrix by replacing the main diagonal by $-\lambda_{j}$, and the lower diagonal by 1 . The singular values of $B$ are easily calculated: since $B^{*} B$ is similar to the tridiagonal Toeplitz matrix with main diagonal $\left|\lambda_{j}\right|^{2}+1$ and lower and upper diagonals $\left|\lambda_{j}\right|$, its eigenvalues are given by $\left|\lambda_{j}\right|^{2}+1-$ $2\left|\lambda_{j}\right| \cos (\pi m /(2 n)), m=1, \ldots, 2 n-1$, and hence

$$
\frac{\|\vec{q}\|}{\|\vec{p}\|}=\frac{\|B \vec{p}\|}{\|\vec{p}\|}\left\{\begin{array}{l}
\leq \sqrt{\left|\lambda_{j}\right|^{2}+1-2\left|\lambda_{j}\right| \cos \left(\pi \frac{2 n-1}{2 n}\right)} \leq 1+\left|\lambda_{j}\right| \\
\geq \sqrt{\left|\lambda_{j}\right|^{2}+1-2\left|\lambda_{j}\right| \cos \left(\frac{\pi}{2 n}\right)}
\end{array} \geq\left(1+\left|\lambda_{j}\right|\right) \sin \left(\frac{\pi}{4 n}\right) \geq \frac{1+\left|\lambda_{j}\right|}{2 n} .\right.
$$

In addition, since $p(z)=\ell_{j}(z)^{2}$, we get from (4) and the Cauchy-Schwarz inequality that

$$
\|\vec{p}\|^{2}=\frac{1}{2 \pi} \int_{|w|=1}|p(w)|^{2}|d w| \geq\left[\frac{1}{2 \pi} \int_{|w|=1}|p(w)||d w|\right]^{2}=\left\|\vec{\ell}_{j}\right\|^{4},
$$

and that

$$
\begin{aligned}
\|\vec{p}\|^{2} & =\frac{1}{2 n} \sum_{m=0}^{2 n-1}\left|p\left(e^{2 \pi i m /(2 n)}\right)\right|^{2} \\
& \leq \frac{1}{2 n}\left(\sum_{m=0}^{2 n-1}\left|p\left(e^{2 \pi i m /(2 n)}\right)\right|\right)^{2}=2 n\left(\frac{1}{2 n} \sum_{m=0}^{2 n-1}\left|\ell_{j}\left(e^{2 \pi i m /(2 n)}\right)\right|^{2}\right)^{2}=2 n\left\|\overrightarrow{\ell_{j}}\right\|^{4} .
\end{aligned}
$$

Consequently, the inequalities (12) are valid, which proves Theorem 2.1.

As a consequence of Theorem 2.1 we can say that, at least up to the factor $\eta_{j, n}$, the illdisposedness of the generalized eigenvalue $\lambda_{j}$ of (1) with respect to either structured and to unstructured perturbations can be measured by checking whether

$$
\begin{equation*}
\rho_{j}:=\left(\left|\lambda_{j}\right|+1\right)\left\|W_{n}^{-1} e_{j}\right\|^{2}\left(\left\|H_{n}\right\|+\left\|\widetilde{H}_{n}\right\|\right) \tag{13}
\end{equation*}
$$

is "large". The term $\eta_{j, n}$ is close to some modest power of $n$ depending on the choice of the norm and on the techniques chosen to prove Theorem 2.1. Though for particular applications it might be desirable to have sharper results, we believe that this factor $\eta_{j, n}$ should be neglected in the measurement of ill-disposedness. As pointed out by Wilkinson [Wil71], "Usually the bound itself is weaker that it might have been because of the necessity of restricting the mass of details to a reasonable level and because of the limitations imposed by expressing the errors in terms of matrix norms."

## 3 Sensitivity of the GEP

In formula (13) of the preceding section we introduced the quantities $\rho_{j}$ for measuring the conditioning of the (nonlinear) generalized eigenvalue problem (1). The aim of this section is to relate our conditioning measures to those used for linear problems, in particular to the condition number of the underlying Hankel matrix $H_{n}$ and the row-scaled Vandermonde matrix $W_{n}$. We are primarily concerned with determining when at least one of the generalized eigenvalues is ill-disposed. Our results will enable us (in the second part of this section) to describe classes of configurations of $\lambda_{1}, \ldots, \lambda_{n}$ where at least one of the quantities $\rho_{j}$ will be large.

### 3.1 Sensitivity of the GEP in terms of condition numbers

The following result is an interesting direct consequence of Theorem 2.1.
Corollary 3.1. If the Hankel matrix $H_{n}$ is ill-conditioned then at least one of the generalized eigenvalues of (1) is ill-disposed. More precisely, we have the estimate

$$
\rho_{1}+\rho_{2}+\ldots+\rho_{n} \geq\left\|H_{n}\right\|\left\|H_{n}^{-1}\right\| .
$$

Proof. From (13) we have that $\rho_{j} \geq\left\|W_{n}^{-1} e_{j}\right\|^{2}\left\|H_{n}\right\|$ for each $j$. Thus

$$
\begin{aligned}
\sum_{j=1}^{n} \rho_{j} & \geq\left\|H_{n}\right\| \sum_{j=1}^{n}\left\|W_{n}^{-1} e_{j}\right\|^{2}=\left\|H_{n}\right\| \operatorname{trace}\left(\left(W_{n}^{-1}\right)^{*} W_{n}^{-1}\right) \\
& \geq\left\|H_{n}\right\|\left\|\left(W_{n}^{-1}\right)^{*} W_{n}^{-1}\right\| \geq\left\|H_{n}\right\|\left\|\left(W_{n}^{t} W_{n}\right)^{-1}\right\|=\left\|H_{n}\right\|\left\|H_{n}^{-1}\right\|
\end{aligned}
$$

Hankel matrices are suspected to be notoriously ill-conditioned (for instance, the famous Hilbert matrix is a Hankel matrix). In the case of positive definite Hankel matrices (or, equivalently, $\lambda_{j} \in \mathbb{R}$ and $c_{j}>0$ for all $j$ ) we can be even more precise. Namely, in [Bec00] it is shown that any positive definite Hankel matrix of order $n$ has a spectral condition number bounded below by $3.210^{n-1} /(16 n)$. Thus Corollary 3.1 confirms a result of Gautschi [Gau82, Gau86] saying that, in case $c_{j}>0$ and $\lambda_{j} \in \mathbb{R}$, the sensitivity of the map $G_{n}$ is increasing exponentially in $n$.

Of course there are also Hankel matrices which are well-conditioned (for example the counter identity), and therefore we should say more about the case of non-real data. Here the following special case occurs quite often in applications.

Definition 3.2. We say that the unit disk case holds if the (unperturbed) moments $h_{k}$ are generated by (3) with $\left|c_{j}\right| \leq 1$ and distinct $\left|\lambda_{j}\right| \leq 1$.

For the case of shape detection from moments we always have $\left|c_{j}\right| \leq 1$ by [GMV99, (2.2)], and the assumption $\left|\lambda_{j}\right| \leq 1$ will be true after scaling of the moments (such a scaling was proposed to be useful in [GMV99, Section 4]). For the application of reconstructing a sparse black box polynomial, we have $\left|\lambda_{j}\right|=1$, and the assumptions on $c_{j}$ can be seen as reasonable assumptions on the scaling of the unknown coefficients $c_{j}$. These applications are discussed later in $\S 4$.

For the unit disk case, we have the following direct consequence of Theorem 2.1.
Corollary 3.3. Suppose that the unit disk case holds. If the row-scaled Vandermonde matrix $W_{n}$ defined in (6) is ill-conditioned then at least one of the generalized eigenvalues of (1) is ill-disposed. More precisely, we have the estimate

$$
\rho_{1}+\ldots+\rho_{n} \geq \frac{\left\|H_{n}\right\|}{n^{2}}\left(\left\|W_{n}\right\|_{F}\left\|W_{n}^{-1}\right\|_{F}\right)^{2}
$$

Proof. From (13) we know that

$$
\rho_{1}+\ldots+\rho_{n} \geq\left\|H_{n}\right\| \sum_{j=1}^{n}\left\|W_{n}^{-1} e_{j}\right\|^{2}=\left\|H_{n}\right\|\left\|W_{n}^{-1}\right\|_{F}^{2}
$$

Thus the assertion follows by observing that in the unit disk case we have the rough bound $\left\|W_{n}\right\|_{F} \leq n$.

The factor $n^{-2}$ in front of the Frobenius condition number of $W_{n}$ in Corollary 3.3 depends on the choice of the norm and the techniques of proof, and thus does not seem to be significant. Also, in typical examples, the quantity $\left\|H_{n}\right\|$ is of magnitude 1 . Thus, roughly speaking, Corollary 3.3 tells us that, within the set of generalized eigenvalues, there is at least one for which the sensitivity is bounded below by the square of the condition number of $W_{n}$.

In comparing Corollaries 3.1 and 3.3 we recall that $H_{n}=W_{n}^{t} W_{n}$. Thus, the square of the Frobenius condition number of $W_{n}$ might potentially be much larger than the condition number of the Hankel matrix $H_{n}$, and so may lead to sharper bounds.

### 3.2 Sensitivity of the GEP with respect to distribution of eigenvalues

In Corollary 3.3 we related the sensitivity of the generalized eigenvalue problem (1) to the Frobenius condition number of the row-scaled Vandermonde matrix $W_{n}$ defined in (6). Clearly, a small $\left|c_{j}\right|$ will deteriorate the condition number of $W_{n}$. However, even in the case of $c_{j}$ minimizing the Frobenius condition number, the condition number may still be "large" depending on the distribution of the abscissa $\lambda_{j}$ in the unit disk. The aim of the remainder of this section is to make this last point more precise.

Definition 3.4. For a given compact set $E$ in the complex plane, define

$$
\begin{equation*}
\gamma_{n}(E)=\inf \left\{\min _{D \text { diagonal }}\left\|D V_{n}\right\|_{F}\left\|\left(D V_{n}\right)^{-1}\right\|_{F}: \lambda_{j} \in E \text { distinct }\right\} . \tag{14}
\end{equation*}
$$

Besides being of interest in its own right, the quantity $\gamma_{n}(E) \geq n$ gives an alternate lower bound for ill-disposedness, as shown in the following result.

Lemma 3.5. Let $E$ be a compact set with $\lambda_{1}, \ldots, \lambda_{n} \in E$. Then

$$
\begin{equation*}
\rho_{1}+\ldots+\rho_{n} \geq \gamma_{n}(E) \tag{15}
\end{equation*}
$$

Moreover, in the unit disk case we have that

$$
\begin{equation*}
\rho_{1}+\ldots+\rho_{n} \geq \frac{\left\|H_{n}\right\|}{n^{2}} \gamma_{n}(E)^{2} \tag{16}
\end{equation*}
$$

Proof. We start by observing that, for fixed $\lambda_{1}, \ldots, \lambda_{n}$, it is not difficult to determine the diagonal factor in (14) which minimizes the Frobenius condition number. Indeed, using the CauchySchwarz inequality we find that

$$
\begin{equation*}
\left\|D V_{n}\right\|_{F}\left\|\left(D V_{n}\right)^{-1}\right\|_{F} \geq \sum_{j=1}^{n}\left\|\left(D V_{n}\right)^{t} e_{j}\right\|\left\|\left(D V_{n}\right)^{-1} e_{j}\right\|=\sum_{j=1}^{n}\left\|V_{n}^{t} e_{j}\right\|\left\|V_{n}^{-1} e_{j}\right\| \tag{17}
\end{equation*}
$$

with equality if $D=\operatorname{diag}\left(\sqrt{\left\|V_{n}^{-1} e_{j}\right\| /\left\|V_{n}^{t} e_{j}\right\|}\right)$.
In order to show (15) we first observe that for the row-scaled Vandermonde matrix $W_{n}$ of (6) we have

$$
\left\|H_{n}\right\|\left\|W_{n}^{-1} e_{j}\right\|=\left\|W_{n}^{t} W_{n}\right\|\left\|W_{n}^{-1} e_{j}\right\| \geq\left\|W_{n}^{t} e_{j}\right\|=\sqrt{\left|c_{j}\right|}\left\|V_{n}^{t} e_{j}\right\|
$$

and $\left\|W_{n}^{-1} e_{j}\right\|=\left\|V_{n}^{-1} e_{j}\right\| / \sqrt{\left|c_{j}\right|}$. Thus, by (13),

$$
\rho_{1}+\ldots+\rho_{n} \geq \sum_{j=1}^{n}\left\|H_{n}\right\|\left\|W_{n}^{-1} e_{j}\right\|^{2} \geq \sum_{j=1}^{n}\left\|V_{n}^{t} e_{j}\right\|\left\|V_{n}^{-1} e_{j}\right\| \geq \gamma_{n}(E)
$$

where the last inequality follows from (17). Finally, in the unit disk case, equation (16) follows directly from Corollary 3.3 and Definition 3.4.

For the remainder of this section we will give lower bounds for $\gamma_{n}(E)$ and thus for the measure of ill-disposedness of our generalized eigenvalues. Notice that if $\partial \mathbb{D} \subset E$, then we may choose as $\lambda_{j}$ the $n$-th roots of unity. In this case $W_{n}=V_{n} / \sqrt{n}$ is unitary, and thus $\gamma_{n}(E)$ attains its smallest possible value, namely $n$. We will show in Theorem 3.6 below that, for sufficiently "nice" sets $E, \gamma_{n}(E)$ can be bounded below and above by a quantity $\gamma(E)^{n}$ multiplied by a modest power of $n$. Here $\gamma(E) \geq 1$ measures which part of the unit circle is not part of $E$. When $E$ contains the unit circle $\partial \mathbb{D}$ then $\gamma(E)=1$. If however $\partial \mathbb{D} \not \subset E$, then $\gamma(E)>1$, and thus $\gamma_{n}(E)$ grows exponentially, with the rate depending on the size of the part of the unit circle not included in $E$.

In order to give a precise definition of $\gamma(E)$, we require some tools from complex analysis. In order to simplify our work we will assume in the following that $E$ is a simply connected compact set. Denote by $\phi$ the Riemann map which maps the complement of $E$ conformally onto the complement of the closed unit disk $\mathbb{D}$. Then the Green function of $E$ with pole at $y$ is given by

$$
g_{E}(z, y)= \begin{cases}\log \left(\left|\frac{1-\overline{\phi(y)} \phi(z)}{\phi(z)-\phi(y)}\right|\right) & \text { for } z, y \in \overline{\mathbb{C}} \backslash E  \tag{18}\\ 0 & \text { otherwise }\end{cases}
$$

with $g_{E}$ being continuous in both variables. Notice that $g_{E}(z, y) \geq 0$ for $z, y \in \mathbb{C}$, and strictly positive if both $z$ and $y$ lie in the complement of $E$. Thus the constant

$$
\begin{equation*}
\gamma(E)=\exp \left(\max _{x \in \partial \mathbb{D}} \frac{1}{2 \pi} \int_{0}^{2 \pi} g_{E}\left(x, e^{i t}\right) d t\right) \tag{19}
\end{equation*}
$$

is at least 1 , with equality if and only if $\partial \mathbb{D} \subset E$. It will be shown in the proof of Lemma 5.2 that in some cases we have a simplified expression for $\gamma(E)$, namely

$$
\begin{equation*}
E \subset \mathbb{D}: \quad \gamma(E)=\exp \left(\max _{x \in \partial \mathbb{D}} g_{E}(x, \infty)\right) \tag{20}
\end{equation*}
$$

For $E$ a real interval, it follows from [Bec00, Remark 3.4] (where the spectral condition number was considered) that

$$
\begin{equation*}
\gamma_{n+1}(E) \geq \frac{1}{2 \sqrt{2 n+2}}\left(\gamma(E)^{n}-2 \gamma(E)^{-n}\right) \tag{21}
\end{equation*}
$$

with the right-hand side clearly being exponentially increasing. Inequality (21) shows, for example, that positive definite Hankel matrices of order $n$ have a condition number bounded below by some term exponentially increasing in $n$. In fact, a behavior similar to (21) is true for any sufficiently "nice" set $E$.

Theorem 3.6. For any compact set $E \subset \mathbb{C}$ which is regular with respect to the Dirichlet problem we have that

$$
\lim _{n \rightarrow \infty} \gamma_{n}(E)^{1 / n}=\gamma(E)
$$

and $\gamma_{n}(E) \leq n^{2} \gamma(E)^{n-1}$ for all $n \geq 1$. If, in addition $E$ is simply connected and of bounded variation $V$, then we have that

$$
\begin{equation*}
\gamma_{n}(E) \geq \frac{1}{\sqrt{n}}\left[\frac{\pi}{2 V} \gamma(E)^{n-1}-1-\frac{V}{\pi}\right] \tag{22}
\end{equation*}
$$

Proof. See Section 5.

With respect to the second part of Theorem 3.6 we notice that any simply connected compact set is known to be regular with respect to the Dirichlet problem. We also recall, for example from [And83], that $E$ is of bounded variation $V$ if, given some parametrization $\beta:[0,1] \mapsto \partial E$ of the boundary of $E$, there exists a tangent at almost every $\beta(s)$, forming an angle $\theta(s)$ with the positive real axis, and if $\theta$ has a total variation bounded by $V$. For instance, convex sets are of bounded variation $2 \pi$.

In order to make the statements of Theorem 3.6 more precise, we need to know both $V$ and the Riemann map. The latter can be determined for various domains, for example for sectors $E=\left\{r e^{i t}: \alpha^{\prime} \leq t \leq \alpha+\alpha^{\prime}, 0 \leq r \leq 1\right\}$ (here $V=2 \pi$ ), for a rectangle $E$ with edges $\pm a \pm i b$ (here $V=2 \pi$ ), or more generally for a polygon (see any advanced textbook on conformal mappings). Here we will have a look at two particular examples.

Example 3.7. Consider the case of a subarc of the unit circle $E=\left\{e^{i t}: \alpha^{\prime} \leq t \leq \alpha+\alpha^{\prime}\right\}$ with $0<\alpha<2 \pi$. Here one easily verifies that $V=2 \pi+2 \alpha$. The Riemann map $\phi$ is found as $\phi_{1} \circ \phi_{2}$, where

$$
\phi_{1}(z)=z / \tan (\alpha / 4)+\sqrt{(z / \tan (\alpha / 4))^{2}-1}
$$

which maps $\overline{\mathbb{C}} \backslash\left[-\tan \left(\frac{\alpha}{4}\right), \tan \left(\frac{\alpha}{4}\right)\right]$ conformally on $\overline{\mathbb{C}} \backslash \mathbb{D}$ and

$$
\phi_{2}(z)=i \frac{e^{-i \alpha^{\prime}-i \alpha / 2} z-1}{e^{-i \alpha^{\prime}-i \alpha / 2} z+1}
$$

which maps $\overline{\mathbb{C}} \backslash E$ conformally on $\overline{\mathbb{C}} \backslash\left[-\tan \left(\frac{\alpha}{4}\right), \tan \left(\frac{\alpha}{4}\right)\right]$. Since $E \subset \mathbb{D}$, we may apply (18) and (20), leading to

$$
\gamma(E)=\max _{\alpha+\alpha^{\prime} \leq t \leq \alpha^{\prime}+2 \pi}\left|\frac{1-\phi\left(e^{i t}\right) \overline{\phi(\infty)}}{\phi\left(e^{i t}\right)-\phi(\infty)}\right|, \quad \phi(\infty)=i \frac{1+\cos \left(\frac{\alpha}{4}\right)}{\sin \left(\frac{\alpha}{4}\right)}, \quad \phi_{1}\left(e^{i\left(t+\alpha^{\prime}+\frac{\alpha}{2}\right)}\right)=i \tan \left(\frac{t}{2}\right) .
$$

Thus

$$
\begin{equation*}
\gamma(E)=|\phi(\infty)|=\frac{1+\cos (\alpha / 4)}{\sin (\alpha / 4)}=\frac{1}{\tan (\alpha / 8)}, \quad 0<\alpha<2 \pi \tag{23}
\end{equation*}
$$

Example 3.8. For an ellipse $E \subset \mathbb{D}$ with foci at $\pm c$ and half axes $c(R \pm 1 / R) / 2, R \geq 1$ (including the case of the interval $[-c, c]$ for $R=1$ ) we have $V=2 \pi$ (since $E$ is convex), and $\phi(z)=\Phi(z / c) / R, \phi(\infty)=\infty$ where $\Phi(z)=z+\sqrt{z^{2}-1}$. If $0<c \leq 2 /(R+1 / R)$ then $E \subset \mathbb{D}$, and a simple computation shows that

$$
\gamma(E)=|\phi(i)|=\frac{1}{R}\left(\frac{1}{c}+\sqrt{\frac{1}{c^{2}}+1}\right)
$$

In particular, for the special case $c=2 /(R+1 / R)$ of ellipses in $\mathbb{D}$ containing $\pm 1$ we find that $\gamma(E)=\left(1+\sqrt{1+c^{2}}\right) /\left(1-\sqrt{1-c^{2}}\right)$. Notice that this is increasing in $c$ (that is as the eccentricity increases), tends to 1 for $c \rightarrow 0$ (that is the case of small eccentricity, $E=\mathbb{D}$ ) and to $1+\sqrt{2}$ for $c \rightarrow 1$ (the case of large eccentricity, $E=[-1,1]$ ). If, on the other hand we have $2 /(R-1 / R) \leq c$, then $\mathbb{D} \subset E$ and hence $\gamma(E)=1$. Finally, if $2 /(R-1 / R)>c>2 /(R+1 / R)$ then we need to apply formula (19), with the maximum being attained for $x=i$ by symmetry. Since the resulting expressions are complicated, we omit details.

We can conclude from (23) and Theorem 3.6 that, if a part of the unit circle is omitted by the generalized eigenvalues of (1), then some of them are necessarily ill-disposed, even if $n$ is not too large. However, the following example shows that for well-disposedness it is not enough to simply require that the generalized eigenvalues are distributed throughout the unit circle and are well separated.


Figure 1: The eigenvalues of Example 3.9 for $m=4$.

Example 3.9. Let $E_{n}=\left\{\lambda_{1}, \ldots, \lambda_{n}\right\}$, with $n=6 m+1$, be the set

$$
\left\{\exp \left(\frac{2 \pi i j}{4 m+1}\right): j=0, \ldots, 4 m\right\} \cup\left\{\exp \left(\frac{2 \pi i(2 j+1)}{8 m+2}\right): j=-m,-m+1, \ldots, m-1\right\},
$$

a subset of the unit circle which has no significant gaps, (see Figure 1 for $m=4$.) We claim that

$$
\begin{equation*}
\gamma_{n}\left(E_{n}\right) \geq \frac{\rho^{m}}{n} \tag{24}
\end{equation*}
$$

for $\rho=\exp \left(\frac{8 \text { Catalan }}{\pi}\right) \approx 10.30>1$, so that there is exponential growth.
To see this claim, suppose without loss of generality that $\lambda_{1}=1$. We make use of (17) and of the formulas (4) and (11) in order to conclude that

$$
\gamma_{n}\left(E_{n}\right) \geq\left\|V_{n}^{t} e_{1}\right\|\left\|V_{n}^{-1} e_{1}\right\|=\sqrt{n}\left\|V_{n}^{-1} e_{1}\right\| \geq \max _{|z|=1}\left|\ell_{1}(z)\right| \geq\left|\ell_{1}(-1)\right|,
$$

with $\ell_{1}$ the correspond first Lagrange polynomial, see (10). With our choice of abscissa we have that

$$
\prod_{j=1}^{n}\left(z-\lambda_{j}\right)=\left(z^{4 m+1}-1\right) \widetilde{\omega}(z), \quad \text { where } \widetilde{\omega}(z)=\prod_{j=-m}^{m-1}\left(z-\exp \left(\frac{2 \pi i(2 j+1)}{8 m+2}\right)\right)
$$

and thus, again by (10),

$$
\left|\ell_{1}(-1)\right|=\frac{|\widetilde{\omega}(-1)|}{(4 m+1)|\widetilde{\omega}(1)|} \geq \frac{|\widetilde{\omega}(-1)|}{n|\widetilde{\omega}(1)|} .
$$

Notice that in the right-hand expression there are no longer the $(4 m+1)$-th roots of unity, and that the zeros of $\widetilde{\omega}$ all lie in the right-hand half circle. This enables us to show exponential growth. Indeed, the mapping $x \mapsto \log (1 / \tan (x))$ is both decreasing and concave on $[0, \pi / 4]$, and
hence

$$
\begin{aligned}
\log \left(\frac{|\widetilde{\omega}(-1)|}{|\widetilde{\omega}(1)|}\right) & =\sum_{j=-m}^{m-1} \log \left(\left|\frac{-1-\exp \left(\frac{2 \pi i(2 j+1)}{8 m+2}\right)}{1-\exp \left(\frac{2 \pi i(2 j+1)}{8 m+2}\right)}\right|\right) \\
& =2 \sum_{j=0}^{m-1} \log \left(\frac{1}{\left|\tan \left(\frac{\pi(2 j+1)}{8 m+2}\right)\right|}\right) \\
& \geq 2 \sum_{j=0}^{m-1} \log \left(\frac{1}{\left|\tan \left(\frac{\pi}{4} \frac{2 j+1}{2 m}\right)\right|}\right) d t \\
& \geq 2 m \int_{0}^{1} \log \left(\frac{1}{\left|\tan \left(\frac{\pi t}{4}\right)\right|}\right) d t=\frac{8 m \text { Catalan }}{\pi} .
\end{aligned}
$$

In the last estimate we have used the classical fact that the midpoint quadrature rule gives an upper bound for the corresponding integral of some concave function. This proves (24). We remark that the term $\rho$ in (24) can in fact be shown to be optimal.

## 4 Applications of the Generalized Hankel Eigenvalue Problem

In this section we briefly describe two applications of the generalized Hankel eigenvalue problem, sparse polynomial interpolation and the shape from moments problem.

### 4.1 Sparse Interpolation of Black Box Polynomials

Sparse interpolation is the problem of finding a sparse standard representation of a multivariate polynomial $f\left(x_{1}, \ldots, x_{d}\right)$ represented as a black-box. That is we want to find constants $c_{i}$ and powers $m_{i, j}$ such that one can write $f$ in its standard form

$$
\begin{equation*}
f\left(x_{1}, \ldots, x_{d}\right)=\sum_{j}^{t} c_{j} x_{1}^{m_{1, j}} \cdots x_{d}^{m_{d, j}} \tag{25}
\end{equation*}
$$

Here the number of nonzero terms $t$ is typically considered as being known. The sparse interplation problem in a numerical context has been introduced and studied by Giesbrecht, Labahn and Lee [GLL04a, GLL06]. Applications of sparse interpolation include its use in approximate polynomial factorization and the decomposition of approximately specified polynomial systems. We refer the reader to the references in [GLL06] for more information on these applications.

Assuming that an upper bound of the degree for each $x_{j}$ is available (which is typically the case), the initial approach found in [GLL06] evaluates the polynomial at the points

$$
h_{k}=f\left(v_{1}^{k}, \ldots, v_{d}^{k}\right), \quad k=0, \ldots, 2 t-1
$$

with each $v_{j}=e^{\frac{2 \pi i}{p_{j}}}$. Here $p_{j}$ are relatively prime integers with $p_{j}$ larger than the upper bound of the degree of term $x_{j}$. If we set $\lambda_{j}=e^{2 \pi i\left(\frac{m_{1, j}}{p_{1}}+\cdots+\frac{m_{d, j}}{p_{d}}\right)}$, then this sparse interpolation problem becomes one of determining the $c_{j}$ and $\lambda_{j}$ satisfying

$$
h_{k}=\sum_{j=1}^{t} c_{j} \lambda_{j}^{k}, \quad k=0, \ldots, 2 t-1
$$

Once the $\lambda_{j}$ are known, the individual powers $m_{i, j}$ are determined by taking logarithms and applying Chinese remaindering. We refer the reader to [GLL06] for further details. Notice that in this application each $\lambda_{j}$ lies on the unit circle.

Example 4.1. Consider the sparse interpolation problem for the simple univariate polynomial

$$
f(x)=\sum_{j=1}^{11} c_{j} x^{j-1}+\sum_{j=1}^{10} c_{11+j} x^{p+j-11}
$$

having 21 terms and $p$ as our upper bound. In this example the corresponding abscissas on the unit circle are $\left\{1, e^{\frac{2 \pi i}{p}}, \ldots, e^{\frac{20 \pi i}{p}}, e^{\frac{-20 \pi i}{p}}, \ldots, e^{\frac{-2 \pi i}{p}}\right\}$. In this case (22) together with Example 3.7 for $\alpha=\frac{20 \pi}{p}$ gives a lower bound for the conditioning of the corresponding generalized eigenvalues. Indeed for a degree as small as $p=101$ the resulting condition number is already bounded below by $10^{20}$, illustrating the conditioning problem for this example.

To reduce the ill-conditioning in the sparse interpolation problem, [GLLO6] adopts a strategy of evaluating points at $v_{j}=e^{2 \pi i \frac{r_{j}}{p_{j}}}$ with $r_{j}$ a random point between 1 and $p_{j}-1$, where $p_{j}$ is a distinct prime and $\lambda_{j}=e^{2 \pi i\left(\frac{m_{1, j} r_{1}}{p_{1}}+\cdots+\frac{m_{d, j} r_{d}}{p_{d}}\right)}$. They show that, for such a random choice of root of unity, the condition number of the associated problem is reasonable with probability at least 1
$\overline{20}$ only ensures a distribution on the half circle.
Here (22) together with Example 3.7 for $\alpha=\pi$ give the lower bound

$$
\gamma_{21}(E) \geq \frac{(1+\sqrt{2})^{20}}{8 \sqrt{21}}-\frac{5}{\sqrt{21}} \approx 1.2310^{6}
$$

which together with (15) and (16) implies that at least some of the $\lambda_{j}$ are quite ill-disposed. The authors in [GLL06] propose to restart the interpolation if such ill-conditioning is encountered, as another random choice will again be well-conditioned with reasonable probability.

### 4.2 The Shape from Moments Problem

The shape from moments problem is to find the vertices $z_{j}$ of a polygonal domain $D$ using complex moments. The problem has a wide variety of uses in both pure and applied mathematic, for example for solving inverse problems for uniform regions related to general elliptical equations. We refer the reader to [GMV99] and the references therein for examples of such applications.

The method takes advantage of a formula originally due to Motzkin and Schoenberg [Dav64, also the references given there] and Davis [Dav64, Dav77] which states that for any simply connected polygonal domain $D$ with $t$ vertices $z_{1}, \ldots, z_{t}$, there are weights $c_{1}, \ldots, c_{t}$ such that for any function $f$ analytic in the closure of $D$

$$
\begin{equation*}
\iint_{D} f^{\prime \prime}(z) d x d y=\sum_{j=1}^{t} c_{j} f\left(z_{j}\right) . \tag{26}
\end{equation*}
$$

In fact, there is an explicit formula for the $c_{j}$ in terms of the $z_{j}$ (see, e.g., [GMV99, Eqn. (2.2)]), and $\left|c_{j}\right|=\left|\sin \left(\psi_{j}\right)\right|$, with $\psi_{j}$ being the outer angle of $D$ at $z_{j}$.

If we let

$$
h_{k}=k(k-1) \iint_{D} z^{k-2} d x d y
$$

be weighted complex moments, then formula (26) becomes

$$
h_{k}=\sum_{j=1}^{t} c_{j} z_{j}^{k}
$$

so that our shape from moments problems attempts to find $c_{1}, \ldots, c_{t}$ and $z_{1}, \ldots, z_{t}$ from the knowledge of $h_{k}$ for $0 \leq k \leq 2 t-1$.

In practice one also transforms these moments via shifts in order to center the problem around the center of mass of the polygonal region. In addition, in [GMV99] these transformed moments are scaled in order that the region can be moved inside the unit disk. If the region is already inside the unit disk then the scaling is done in such a way as to enlarge the region so as to take more space inside the disk.

As a consequence of the scaling used in the shape from moments problem the unit disk case holds, and so implies that $\left\|H_{n}\right\|+\left\|\widetilde{H}_{n}\right\| \leq 2 n^{2}$. In addition, $\left\|H_{n}\right\| \geq\left|h_{2}\right|$, which by (26) equals twice the volume of $D$. Thus, in view of (13), the quantity

$$
\left\|W_{n}^{-1} e_{j}\right\|^{2}=\frac{1}{\left|\sin \left(\psi_{j}\right)\right|}\left\|V_{n}^{-1} e_{j}\right\|^{2}
$$

is a good measure for ill-disposedness of the generalized eigenvalue $\lambda_{j}$ of problem (1), or, in other words, of the sensitivity of the problem of recovering the vertex $\lambda_{j}$ from the moments. The following example shows that, even for domains where $\left|\sin \left(\psi_{j}\right)\right|$ is not too small, the quantities $\left\|V_{n}^{-1} e_{j}\right\|$ are quite often exponentially increasing in $n$, implying that the corresponding vertices are ill-disposed.

Example 4.2. Let $\left\{\lambda_{1}, \ldots, \lambda_{n}\right\}$ with $n=2 m$ be the set

$$
\lambda_{2 j}=R e^{2 \pi i \frac{j}{m}} \quad \text { and } \quad \lambda_{2 j-1}=r e^{\pi i \frac{(2 j-1)}{m}} \quad j=1, \ldots, m
$$

with $1 \geq R \geq r>0$. Consider the shape from moments problem used to reconstruct the star pattern with $m$ equally spaced vertices $\lambda_{2 j-1}$ on an inside circle of radius $r$ and $m$ equally spaced vertices $\lambda_{2 j}$ on an outside circle of radius $R$. See Figure 2 for $n=18$.

In order to describe the sensitivity of the problem of recovering vertex $\lambda_{j}$ from the moments, we determine the asymptotic behavior of $\left\|V_{n}^{-1} e_{2 j}\right\|$ and of $\left\|V_{n}^{-1} e_{2 j-1}\right\|$ (which by symmetry do not depend on j) for $n \rightarrow \infty$. We again make use of the formulas (10) and (11) from Section 2 expressing $V_{n}^{-1} e_{j}$ in terms of Lagrange polynomials. Here we observe that the polynomial $\omega$ of (10) can be factored as $\omega(z)=\omega_{1}(z) \omega_{2}(z)$ where

$$
\omega_{1}(z)=\prod_{j=1}^{m}\left(z-\lambda_{2 j-1}\right)=z^{m}+r^{m}, \quad \omega_{2}(z)=\prod_{j=1}^{m}\left(z-\lambda_{2 j}\right)=z^{m}-R^{m}
$$

Consequently, $\left|\omega^{\prime}\left(\lambda_{2 j}\right)\right|=m R^{m-1}\left(R^{m}+r^{m}\right)$, and, by (11),

$$
\left\|V_{n}^{-1} e_{2 j}\right\|^{2}=\frac{\left\|\omega(z) /\left(z-\lambda_{2 j}\right)\right\|_{L_{2}(\partial \mathbb{D})}^{2}}{\left|\omega^{\prime}\left(\lambda_{2 j}\right)\right|^{2}}=\frac{\left\|\omega_{1}(z) \sum_{k=0}^{m-1} z^{m-1-k} \lambda_{2 j}^{k}\right\|_{L_{2}(\partial \mathbb{D})}^{2}}{\left|\omega^{\prime}\left(\lambda_{2 j}\right)\right|^{2}}=\frac{\epsilon_{m}}{m^{2} R^{4 m-2}}
$$

where $\epsilon_{m} \in[1 / 4,2 m]$. Similarly one can show that

$$
\left\|V_{n}^{-1} e_{2 j-1}\right\|^{2}=\frac{\left\|\omega_{2}(z) \sum_{k=0}^{m-1} z^{m-1-k} \lambda_{2 j-1}^{k}\right\|_{L_{2}(\partial \mathbb{D})}^{2}}{\left|m r^{m-1}\left(r^{m}+R^{m}\right)\right|^{2}}=\frac{\epsilon_{m}^{\prime}}{m^{2} r^{2 m-2} R^{2 m}}
$$



Figure 2: The polygon of Example 4.2 for $n=18$ with $R=1$ and $r=\frac{1}{}$

## 5 Proof for Theorem 3.6

Our proof of Theorem 3.6 is divided in several parts. First, in Lemma 5.1 we relate the quantity $\gamma_{n}(E)$ defined in (14) to the solution $\tilde{\gamma}_{n}(E)$ of an extremal problem for weighted polynomials in the complex plane, namely, the problem of finding the quantity

$$
\begin{equation*}
\tilde{\gamma}_{n}(E):=\max \left\{\frac{\|P\|_{L_{\infty}(\partial \mathbb{D})}}{\left\|\rho^{n} P\right\|_{L_{\infty}(E)}}: P \text { is a polynomial of degree } \leq n\right\}, \tag{27}
\end{equation*}
$$

where $\rho(z):=1 / \max (|z|, 1)$. In Lemma 5.2 we report about some recent results on weighted polynomials and relate the quantity $\tilde{\gamma}_{n}(E)$ to $\gamma(E)$ defined in (19), (20). For the special case of simply connected $E \subset \mathbb{D}$ we give in Lemma 5.3 sharper estimates for a quantity closely related to $\tilde{\gamma}_{n}(E)$. Finally, we follow the Lemma with a proof of Theorem 3.6.

Lemma 5.1. Let $E \subset \mathbb{C}$ be compact and $n \geq 1$ an integer. Then

$$
\frac{1}{\sqrt{n}} \tilde{\gamma}_{n-1}(E) \leq \gamma_{n}(E) \leq n^{2} \tilde{\gamma}_{n-1}(E)
$$

Proof. Consider $\lambda_{1}, \ldots, \lambda_{n} \in E$, and the corresponding Lagrange polynomials $\ell_{j}$ as defined in (10). For any polynomial $P$ of degree at most $n-1$ we obtain using the Lagrange interpolation
formula and (4)

$$
\begin{aligned}
\|P\|_{L_{\infty}(\partial \mathbb{D})} \leq \sqrt{n}\|P\|_{L_{2}(\partial \mathbb{D})} & \leq \sqrt{n} \sum_{j=1}^{n}\left|P\left(\lambda_{j}\right)\right|\left\|\ell_{j}\right\|_{L_{2}(\partial \mathbb{D})} \\
& \leq \sqrt{n}\left\|\rho^{n-1} P\right\|_{L_{\infty}(E)} \sum_{j=1}^{n} \frac{1}{\rho\left(\lambda_{j}\right)^{n-1}}\left\|\ell_{j}\right\|_{L_{2}(\partial \mathbb{D})}
\end{aligned}
$$

Since $\left\|\ell_{j}\right\|_{L_{2}(\partial \mathbb{D})}=\left\|V_{n}^{-1} e_{j}\right\|$ by (11) and $\left\|V_{n}^{t} e_{j}\right\|=\left[1+\left|\lambda_{j}\right|^{2}+\ldots+\left|\lambda_{j}\right|^{2 n-2}\right]^{1 / 2} \geq 1 / \rho\left(\lambda_{j}\right)^{n-1}$ with the weight function $\rho(z)=1 / \max (1,|z|)$ as defined above, it follows that

$$
\|P\|_{L_{\infty}(\partial \mathbb{D})} \leq \sqrt{n}\left\|\rho^{n-1} P\right\|_{L_{\infty}(E)} \gamma_{n}(E)
$$

for any polynomial $P$ of degree at most $n-1$. Thus $\tilde{\gamma}_{n-1}(E) \leq \sqrt{n} \gamma_{n}(E)$, as claimed in the assertion of Lemma 5.1.

In order to show the other estimate, we first have to choose "good" points $\lambda_{j} \in E$ such that the discrete set $E_{n}:=\left\{\lambda_{1}, \ldots, \lambda_{n}\right\}$ satisfies $\tilde{\gamma}_{n-1}\left(E_{n}\right) \approx \tilde{\gamma}_{n-1}(E)$. Here we will take the set of weighted Fekete points ${ }^{4}$ defined as follows (cf. [SaTo97, Section III.1]). If we explicitly denote $V_{n}\left(\lambda_{1}, \ldots, \lambda_{n}\right)$ for a Vandermonde matrix with abscissa $\lambda_{1}, \ldots, \lambda_{n}$, then the weighted Fekete points are those points in $E$ maximizing the expression

$$
\operatorname{det} V_{n}\left(\lambda_{1}, \ldots, \lambda_{n}\right) \prod_{k=1}^{n} \rho\left(\lambda_{k}\right)^{n-1}
$$

For the corresponding Lagrange polynomials $\ell_{j}$ we have

$$
\left\|\rho^{n-1} \ell_{j}\right\|_{L_{\infty}(E)}=\rho^{n-1}\left(\lambda_{j}\right) \max _{z \in E} \frac{\rho(z)^{n-1}}{\rho\left(\lambda_{j}\right)^{n-1}} \frac{\operatorname{det} V_{n}\left(\lambda_{1}, \ldots, \lambda_{j-1}, z, \lambda_{j+1}, \ldots, \lambda_{n}\right)}{\operatorname{det} V_{n}\left(\lambda_{1}, \ldots, \lambda_{n}\right)} \leq \rho^{n-1}\left(\lambda_{j}\right)
$$

by definition of the Fekete points. Hence, for any polynomial $P$ of degree at most $n-1$, by the Lagrange interpolation formula,

$$
1 \leq \frac{\left\|\rho^{n-1} P\right\|_{L_{\infty}(E)}}{\left\|\rho^{n-1} P\right\|_{L_{\infty}\left(E_{n}\right)}} \leq \sum_{j=1}^{n} \frac{\left|P\left(\lambda_{j}\right)\right|}{\left\|\rho^{n-1} P\right\|_{L_{\infty}\left(E_{n}\right)}}\left\|\rho^{n-1} \ell_{j}\right\|_{L_{\infty}(E)} \leq n
$$

implying that $\tilde{\gamma}_{n-1}(E) \leq \tilde{\gamma}_{n-1}\left(E_{n}\right) \leq n \tilde{\gamma}_{n-1}(E)$. Also, by (the sharpness of) the Hölder inequality applied to the Lagrange interpolation formula, we have that, for any $\zeta \in \mathbb{C}$,

$$
\max \left\{\frac{|P(\zeta)|}{\left\|\rho^{n} P\right\|_{L_{\infty}\left(E_{n}\right)}}: P \text { is a polynomial of degree } \leq n-1\right\}=\sum_{j=1}^{\infty}\left|\ell_{j}(\zeta)\right| / \rho\left(\lambda_{j}\right)^{n-1}
$$

Thus

$$
\begin{aligned}
n \tilde{\gamma}_{n-1}(E) & \geq \tilde{\gamma}_{n-1}\left(E_{n}\right)=\max _{\zeta \in \partial \mathbb{D}} \sum_{j=1}^{n}\left|\ell_{j}(\zeta)\right| / \rho\left(\lambda_{j}\right)^{n-1} \\
& \geq\left[\max _{\zeta \in \partial \mathbb{D}} \sum_{j=1}^{n}\left|\ell_{j}(\zeta)\right|^{2} / \rho\left(\lambda_{j}\right)^{2 n-2}\right]^{1 / 2} \\
& \geq\left[\sum_{j=1}^{n} \frac{1}{\rho\left(\lambda_{j}\right)^{2 n-2}} \frac{1}{2 \pi} \int_{|\zeta|=1}\left|\ell_{j}(\zeta)\right|^{2}|d \zeta|\right]^{1 / 2}=\left[\sum_{j=1}^{n} \frac{\left\|\vec{\ell}_{j}\right\|^{2}}{\rho\left(\lambda_{j}\right)^{2 n-2}}\right]^{1 / 2} \\
& \geq \frac{1}{\sqrt{n}} \sum_{j=1}^{n} \frac{\left\|\vec{\ell}_{j}\right\|}{\rho\left(\lambda_{j}\right)^{n-1}} .
\end{aligned}
$$

[^1]On the other hand, by (14), (17), and (11),

$$
\gamma_{n}(E) \leq \gamma_{n}\left(E_{n}\right)=\sum_{j=1}^{n}\left\|\vec{\ell}_{j}\right\| \sqrt{1+\left|\lambda_{j}\right|^{2}+\ldots+\left|\lambda_{j}\right|^{2 n-2}} \leq \sqrt{n} \sum_{j=1}^{n} \frac{\left\|\vec{\ell}_{j}\right\|}{\rho\left(\lambda_{j}\right)^{n-1}}
$$

and consequently the second inequality of Lemma 5.1 is true.

The following result on the weighted polynomial extremal problem (27) was obtained in a more general setting by Maskhar and Saff in 1985. We refer the reader to [SaTo97, Section III.2] and the references therein.

Lemma 5.2. Let $E \subset \mathbb{C}$ be compact, and regular with respect to the Dirichlet problem. Then for any integer $n \geq 1$ we have that $\tilde{\gamma}_{n}(E) \leq \gamma(E)^{n}$, and

$$
\lim _{n \rightarrow \infty} \tilde{\gamma}_{n}(E)^{1 / n}=\gamma(E)
$$

Proof. The logarithmic potential of a probability measure $\mu$ is defined by

$$
U^{\mu}(z):=\int \log \left(\frac{1}{|z-t|}\right) d \mu(t)
$$

According to [SaTo97, Theorem I.1.3 and Theorem I.4.8], given some function $Q$ continuous on $E$ there exists a unique probability measure $\mu_{Q}$ (called the equilibrium measure in the presence of an external field $Q$ ) supported on $E$ and a constant $c_{Q}$ such that

$$
U^{\mu_{Q}}(z)+Q(z) \begin{cases}=c_{Q} & \text { for } z \text { lying in the support of } \mu_{Q} \\ \geq c_{Q} & \text { for } z \in E\end{cases}
$$

In addition, $\mu_{Q}$ has a continuous potential by the assumption on $E$. The weighted BernsteinWalsh inequality [SaTo97, Theorem III.2.1] tells us that, for any integer $n \geq 0$, for any polynomial $P_{n}$ of degree at most $n$, and for any $z \in \mathbb{C}$

$$
\left|P_{n}(z)\right| \leq e^{n\left(c_{Q}-U^{\mu} Q(z)\right)}\left\|e^{-n Q} P_{n}\right\|_{L_{\infty}(E)}
$$

Moreover, this inequality is sharp since by [SaTo97, Theorem III.5.3] and the following remarks we have for all $z \in \mathbb{C}$

$$
e^{c_{Q}-U^{\mu} Q}(z)=\lim _{n \rightarrow \infty} \sup \left\{\left[\frac{\left|P_{n}(z)\right|}{\left\|e^{-n Q} P_{n}\right\|_{L_{\infty}(E)}}\right]^{1 / n}: P_{n} \text { a polynomial of degree } n\right\}
$$

Thus the assertion of the Lemma 5.2 follows by showing that the constant $\gamma(E)$ of (19) equals

$$
\begin{equation*}
\gamma(E)=\max _{z \in \partial \mathbb{D}} \exp \left(c_{Q}-U^{\mu_{Q}}(z)\right) \tag{28}
\end{equation*}
$$

for our external field $Q=\log (1 / \rho)$. Using [SaTo97, Example 0.5.7] we have that

$$
Q(z)=\log \left(\frac{1}{\rho(z)}\right)=\log (\max (1,|z|))=-\int_{-\pi}^{\pi} \log \left(\frac{1}{|z-t|}\right) \frac{d t}{2 \pi}
$$

that is, the external field coincides with $(-1)$ times the logarithmic potential $U^{\nu}$ of a probability measure $\nu$. Thus we know from [SaTo97, Example II.4.8]) that $\mu_{Q}=\hat{\nu}$, the balayage measure of $\nu$ onto $E$ (see, e.g., [SaTo97, Theorem II.4.7]), and

$$
\begin{aligned}
& U^{\hat{\nu}}(z)+Q(z)=U^{\hat{\nu}}(z)-U^{\nu}(z)=F(\infty)-F(z), \quad \text { where } \\
& F(z):=\int g_{E}(z, t) d \nu(t)=\frac{1}{2 \pi} \int_{0}^{2 \pi} g_{E}\left(z, e^{i t}\right) d t
\end{aligned}
$$

The latter formula follows from integration [SaTo97, Equation (4.32) of Chapter II] with respect to $\nu$. Notice that, since the Green function is continuous, the function $F$ equals zero on $E$, and hence $c_{Q}=F(\infty)$. Finally, by equation (19) we have that $\gamma(E)=\|F\|_{L_{\infty}(\partial \mathbb{D})}$, and $Q$ is identically zero on $\partial \mathbb{D}$. Thus (28) holds.

It only remains to establish formula (20). Let $E \subset \mathbb{D}$. Then $U^{\nu}$ is constant on $E$, and hence $\hat{\nu}$ is the equilibrium measure of $E$. Using the classical link between Green functions and potentials of equilibrium measures (cf. [SaTo97, Equation (4.8) of Chapter I]), we conclude that $F(z)=g_{\partial \mathbb{D}}(z, \infty)-g_{E}(z, \infty)+c$ for some constant $c$. Comparing the values on $E$ yields that $c=0$ and so implies the alternate representation (20) in the case $E \subset \mathbb{D}$.

Lemma 5.3. Let $E \subset \mathbb{C}$ be a compact simply connected set of bounded rotation $V$. As before, let $\phi$ map $\overline{\mathbb{C}} \backslash E$ conformally to $\overline{\mathbb{C}} \backslash \mathbb{D}$, with $\phi(\infty)=\infty$, and $\phi^{\prime}(\infty)>0$, and let

$$
Q(x)=\prod_{k=1}^{M}\left(x-\omega_{k}\right), \quad \omega_{k} \in \mathbb{C} \backslash E
$$

Then for all $z \notin E$ and $m \geq M$

$$
\max \left\{\frac{|P(z) / Q(z)|}{\|P / Q\|_{L_{\infty}(E)}}: \operatorname{deg} P \leq m\right\} \geq \frac{\pi}{V}|\phi(z)|^{m-M} \prod_{j=1}^{M}\left|\frac{1-\phi(z) \overline{\phi\left(\omega_{k}\right)}}{\phi(z)-\phi\left(\omega_{k}\right)}\right|-1-\frac{\pi}{V}
$$

Proof. Consider the function

$$
\tilde{R}(w)=w^{m-M} \prod_{k=1}^{M} \frac{1-w \overline{\phi\left(\omega_{k}\right)}}{w-\phi\left(\omega_{k}\right)}
$$

which is a Blaschke product, analytic in $\mathbb{D}$, and of modulus 1 on $\partial \mathbb{D}$, and so in particular satisfies

$$
\|\tilde{R}\|_{L_{\infty}(\partial \mathbb{D})}=1
$$

Following [ElSa83, And83] we also consider the Faber map $T$ which for a rational function $R$ analytic in $|w| \leq \theta, \theta>1$, is defined by

$$
r(z):=T(R)(z)=\frac{1}{2 \pi i} \int_{|\phi(\zeta)|=\theta} \frac{R(\phi(\zeta))}{\zeta-z} d \zeta, \quad|\phi(z)|<\theta
$$

It is known (see, e.g., [ElSa83, Lemma 1.1 and Proposition 1.2]) that, with $R$, also $T(R)$ is rational, and $z \mapsto T(R)(z)-R(\phi(z))$ has an analytic continuation outside $E$, including infinity. In particular, $T(R)$ has all its poles outside of $E$, namely at the images under $\psi$, the inverse mapping of $\phi$, of $R$. For the operator norm we know from [And83] that

$$
\frac{\|T(R)\|_{L_{\infty}(\partial E)}}{\|R\|_{L_{\infty}(\partial \mathbb{D})}} \leq\|T\| \leq \frac{V}{\pi}
$$

As a consequence, the function $\tilde{P}$ defined by

$$
\tilde{P} / Q=\tilde{r}=T(\tilde{R})
$$

is a polynomial of degree at most $m$, and, with the help of the maximum principle, we get for $z \notin E$

$$
\begin{aligned}
\frac{|\tilde{r}(z)|}{\|\tilde{r}\|_{L_{\infty}(E)}} & \geq \frac{\pi}{V} \frac{|\tilde{R}(\phi(z))|-|\tilde{r}(z)-\tilde{R}(\phi(z))|}{\|\tilde{R}\|_{L_{\infty}(\partial \mathbb{D})}}=\frac{\pi}{V}[|\tilde{R}(\phi(z))|-|\tilde{r}(z)-\tilde{R}(\phi(z))|] \\
& \geq \frac{\pi}{V}\left[|\tilde{R}(\phi(z))|-\|\tilde{r}-\tilde{R} \circ \phi\|_{L_{\infty}(\partial E)}\right] \\
& \geq \frac{\pi}{V}\left[|\tilde{R}(\phi(z))|-\|\tilde{r}\|_{L_{\infty}(\partial E)}-\|\tilde{R}\|_{L_{\infty}(\partial \mathbb{D})}\right] \\
& \geq \frac{\pi}{V}|\tilde{R}(\phi(z))|-1-\frac{\pi}{V},
\end{aligned}
$$

as claimed in Lemma 5.3.

Proof of Theorem 3.6 The first part of the Theorem follows by combining Lemma 5.1 and Lemma 5.2. Suppose now that $E$ is simply connected and of bounded variation. We consider first the case $\partial \mathbb{D} \subset E$. Since $\left\|V_{n}^{-1} e_{j}\right\|\left\|V_{n}^{t} e_{j}\right\| \geq\left|e_{j}^{t} V_{n} V_{n}^{-1} e_{j}\right|=1$, we see from (17) that $\gamma_{n}(E) \geq n$. On the other hand, for the abscissa $\lambda_{1}, \ldots, \lambda_{n}$ being the $n$-th roots of unity we know that $V_{n}^{-1}=\frac{1}{n} V_{n}^{*}$ and $\left\|V_{n}^{*} e_{j}\right\|=\sqrt{n}$, and thus $\gamma_{n}(E)=n$ by (14) and (17). Thus in this case the last part of Theorem 3.6 is trivially true since $\gamma(E)=1$.

If $\partial \mathbb{D} \not \subset E$ then $\gamma(E)>1$ by (19), and there exists a $z^{\prime} \in \partial \mathbb{D} \backslash E$ with

$$
\log (\gamma(E))=\frac{1}{2 \pi} \int_{0}^{2 \pi} g_{E}\left(z^{\prime}, e^{i t}\right) d t=\frac{1}{2 \pi} \int_{0}^{2 \pi} g\left(e^{i t}\right) d t
$$

where $g(z):=g_{E}\left(z^{\prime}, z\right)-\log \left(\frac{1}{\left|z^{\prime}-z\right|}\right)$, a continuous function on $\partial \mathbb{D}$. Notice that

$$
\log \left(\gamma(E)^{n}\right)=\frac{1}{\int_{0}^{2 \pi / n} d t} \int_{0}^{2 \pi / n} \sum_{k=1}^{n} g\left(e^{2 \pi i / k} e^{i t}\right) d t \leq \sum_{k=1}^{n} g\left(e^{2 \pi i / k} e^{i t^{\prime}}\right)
$$

for some $t^{\prime} \in[0,2 \pi / n]$. Since $g$ is continuous we may suppose without loss of generality that $z^{\prime} \notin\left\{e^{2 \pi i / k} e^{i t^{\prime}}: k=1, \ldots, n\right\}$. Taking into account (18) and writing $\omega_{k}=e^{2 \pi i / k} e^{i t^{\prime}}$, we conclude that

$$
\begin{equation*}
\gamma(E)^{n} \leq\left|\left(z^{\prime}\right)^{n}-e^{i t^{\prime} n}\right| \prod_{\omega_{k} \notin E}\left|\frac{1-\phi\left(z^{\prime}\right) \overline{\phi\left(\omega_{k}\right)}}{\phi\left(z^{\prime}\right)-\phi\left(\omega_{k}\right)}\right| \tag{29}
\end{equation*}
$$

Let

$$
Q(z)=\prod_{\omega_{k} \notin E}\left(z-\omega_{k}\right), \quad \tilde{Q}(z)=\prod_{\omega_{k} \in E}\left(z-\omega_{k}\right), \quad m:=\operatorname{deg} Q
$$

and observe that, for $z \in E$,

$$
\rho(z)^{n}|Q(z)||\tilde{Q}(z)|=\frac{\left|z^{n}-e^{i t^{\prime} n}\right|}{\max \left\{1,|z|^{n}\right\}} \leq 2
$$

Hence,

$$
\begin{aligned}
\gamma_{n}(E) & \geq \max \left\{\frac{\left|P\left(z^{\prime}\right)\right|}{\left\|\rho^{n} P\right\|_{L_{\infty}(E)}}: \operatorname{deg} P \leq n\right\} \\
& \geq \max \left\{\frac{\left|\tilde{Q}\left(z^{\prime}\right) P\left(z^{\prime}\right)\right|}{\left\|\rho^{n} \tilde{Q} P\right\|_{L_{\infty}(E)}}: \operatorname{deg} P \leq m\right\} \\
& \geq \frac{\left|\left(z^{\prime}\right)^{n}-e^{i t^{\prime} n}\right|}{2} \max \left\{\frac{\left|P\left(z^{\prime}\right) / Q\left(z^{\prime}\right)\right|}{\|P / Q\|_{L_{\infty}(E)}}: \operatorname{deg} P \leq m\right\}
\end{aligned}
$$

Replacing the right-hand expression by the lower bound found in Lemma 5.3 for $M=m$ and using (29) together with the trivial estimate $\left|\left(z^{\prime}\right)^{n}-e^{i t^{\prime} n}\right| \leq 2$, we obtain

$$
\frac{\pi}{2 V} \gamma(E)^{n}-1-\frac{\pi}{V} \leq \gamma_{n}(E) \leq \gamma(E)^{n}
$$

the right-hand bound having been already established in Lemma 5.2. Together with Lemma 5.1 we may conclude that the second part of Theorem 3.6 is also true in the case $\partial \mathbb{D} \not \subset E$.

## 6 Conclusions

In this paper we have shown that structured and unstructured errors of the generalized eigenvalue problem for Hankel matrices (1) are bounded by essentially the same quantities. We also relate the relative sensitivity of the most sensitive eigenvalue to the condition number of the underlying Hankel matrix. Finally, a new result on the growth of the condition number of special scaled Vandermonde matrices with complex nodes is also given. This result shows that the conditioning of the generalized eigenvalue problem (1) for a subclass of important problems of interest grows exponentially in terms of a quantity that depends on the largest distance between eigenvalues.

In this paper we have analyzed only the sensitivity with respect to ordinary moments. It is natural to ask what can be done with modified moments. Recall that a well-known procedure of improving the sensitivity of the generalized eigenvalue problem is to consider the equivalent problem

$$
T_{n}^{t}\left(\widetilde{H}_{n}-\lambda H_{n}\right) T_{n} y=0
$$

where $T_{n}$ is some suitably chosen invertible upper triangular matrix. By considering the ( $k+1$ )th column of $T_{n}$ as the coefficient vector of some polynomial $p_{k}$ of degree $k$ for $k=0, \ldots, n-1$, we have from (3) that this new generalized eigenvalue problem can be rewritten as $(\widetilde{\mathbf{H}}-\lambda \mathbf{H}) y=0$, where the entries of $\widetilde{\mathbf{H}}=T_{n}^{t} \widetilde{H}_{n} T_{n}$, and $\mathbf{H}=T_{n}^{t} H_{n} T_{n}$, respectively, are given by

$$
\widetilde{\mathbf{H}}_{k, \ell}=\sum_{j=1}^{n} c_{j} \lambda_{j} p_{k}\left(\lambda_{j}\right) p_{\ell}\left(\lambda_{j}\right), \quad \mathbf{H}_{k, \ell}=\sum_{j=1}^{n} c_{j} p_{k}\left(\lambda_{j}\right) p_{\ell}\left(\lambda_{j}\right) .
$$

Such matrices are usually referred to as modified Gramians. The entries $\sum_{j=1}^{n} c_{j} p_{k}\left(\lambda_{j}\right)$ of the (possibly scaled) first column/row of $\mathbf{H}$ are usually referred to as modified moments.

Following the ideas of $\S 2.2$ one may show that an (unstructured) $\epsilon$-perturbation of these two matrices similar to (7) leads to a perturbation of the $j$-th eigenvalue, which in first order can be bounded above by

$$
\left|\frac{d \lambda_{j}}{d \epsilon}(0)\right| \leq\left(\left|\lambda_{j}\right|+1\right)\left\|\mathbf{W}^{-1} e_{j}\right\|^{2}, \quad \mathbf{W}=\operatorname{diag}\left(\sqrt{c_{j}}\right) \mathbf{V}, \quad \mathbf{V}=\left(p_{k-1}\left(\lambda_{j}\right)\right)_{j, k=1, \ldots, n}
$$

That is, we have to replace the Vandermonde matrix in (9) by a generalized Vandermonde matrix. The norm of the inverse and the condition number of such (possibly row-scaled) generalized Vandermonde matrices has been discussed by a number of authors for various choices of bases and particular abscissa (for a summary see for instance [Bec96]). A good choice of a family of polynomials $p_{k}$ may dramatically improve the sensitivity of the generalized eigenvalue problem. For instance, we can produce a unitary $\mathbf{W}$ by taking as $p_{k}$ the orthonormal polynomials with respect to the hermitian scalar product $\ll p, q \gg=\sum_{j=1}^{n}\left|c_{j}\right| \overline{p\left(\lambda_{j}\right)} q\left(\lambda_{j}\right)$. As another example, choosing the corresponding formal orthogonal polynomials leads to $\mathbf{H}$ being the identity matrix (and thus $\mathbf{W}^{-1}=\mathbf{W}^{t}$ ) and $\widetilde{\mathbf{H}}$ being a (in general non hermitian) tridiagonal matrix. Here

$$
\left\|\mathbf{W}^{-1} e_{j}\right\|^{2}=\frac{\sum_{\ell=0}^{n-1}\left|p_{\ell}\left(\lambda_{j}\right)\right|^{2}}{\left|\sum_{\ell=0}^{n-1} p_{\ell}\left(\lambda_{j}\right)^{2}\right|}
$$

and thus the sensitivity will depend on the growth behaviour of the formal orthogonal polynomials of the support of (formal) orthogonality. This in turn is a subject of research where (up to the simple case of real orthogonality $\lambda_{j} \in \mathbb{R}$ and $c_{j}>0$ ) very few results are known.

In summary, a suitable choice of the family of polynomials may lead to a sparse or structured eigenvalue problem or to well-conditioned generalized eigenvalue problems (or sometimes even both). Unfortunately, in general for our applications the entries of the modified Gramians are not available and need to be computed which is a potential new source of errors. In addition, it seems to us that in problems like shape detection or sparse interpolation in general we do not have sufficient a priori information to design a "good" family of polynomials $p_{k}$. As such it remains a topic for future research to alter our generalized eigenvalue problem in order to arrive at better behaved numerical procedures for effective computation.

Of course there are other generalized eigenvalue problems that remain of interest to us for specific applications. As an example, sparse interpolation problems in computer algebra often look for other bases in which sparse representations may be possible. These include bases in terms of Cheyshev polynomials, Pochhammer (factorial) polynomials and others (c.f. [GLL04b]). It would be of interest to study the structured perturbations of these problems and the geometric properties of the set of generalized eigenvalues that might effect condition numbers of such problems.

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[^1]:    ${ }^{4}$ For special sets $E$ some other points like Fejer points or suitable alternants might lead to sharper estimates. Compare with [Bec00].

