

A Riemann–Hilbert Approach to Asymptotic Analysis of Toeplitz+Hankel Determinants

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Abstract. In this paper we will formulate 4×4 Riemann–Hilbert problems for Toeplitz+Hankel determinants and the associated system of orthogonal polynomials, when the Hankel symbol is supported on the unit circle and also when it is supported on an interval $[a, b]$, $0 < a < b < 1$. The distinguishing feature of this work is that in the formulation of the Riemann–Hilbert problem no specific relationship is assumed between the Toeplitz and Hankel symbols. We will develop nonlinear steepest descent methods for analysing these problems in the case where the symbols are smooth (i.e., in the absence of Fisher–Hartwig singularities) and admit an analytic continuation in a neighborhood of the unit circle (if the symbol’s support is the unit circle). We will finally introduce a model problem and will present its solution requiring certain conditions on the ratio of Hankel and Toeplitz symbols. This in turn will allow us to find the asymptotics of the norms h_n of the corresponding orthogonal polynomials and, in fact, the large n asymptotics of the polynomials themselves. We will explain how this solvable case is related to the recent operator-theoretic approach in [Basor E., Ehrhardt T., in Large Truncated Toeplitz Matrices, Toeplitz Operators, and Related Topics, *Oper. Theory Adv. Appl.*, Vol. 259, Birkhäuser/Springer, Cham, 2017, 125–154, arXiv:1603.00506] to Toeplitz+Hankel determinants. At the end we will discuss the prospects of future work and outline several technical, as well as conceptual, issues which we are going to address next within the 4×4 Riemann–Hilbert framework introduced in this paper.

Key words: Toeplitz+Hankel determinants; Riemann–Hilbert problem; asymptotic analysis

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1 Introduction and preliminaries

The $n \times n$ Toeplitz and Hankel matrices associated respectively to the symbols ϕ and w , supported on the unit circle \mathbb{T} are respectively defined as

$$T_n[\phi; r] := \{\phi_{j-k+r}\}, \quad j, k = 0, \dots, n-1, \quad \phi_k = \int_{\mathbb{T}} z^{-k} \phi(z) \frac{dz}{2\pi i z},$$

and

$$H_n[w; s] := \{w_{j+k+s}\}, \quad j, k = 0, \dots, n-1, \quad w_k = \int_{\mathbb{T}} z^{-k} w(z) \frac{dz}{2\pi i z}, \quad (1.1)$$

for fixed *offset* values $r, s \in \mathbb{Z}$. If the Hankel symbol w is supported on a subset I of the real line, then w_k in (1.1) are instead given by

$$w_k = \int_I x^k w(x) dx. \quad (1.2)$$

The Toeplitz and Hankel determinants characterize important objects particularly in random matrix theory, statistical mechanics, theory of orthogonal polynomials, theory of Fredholm determinants, etc. For more on the history of the development of the theory of Toeplitz and Hankel determinants and their numerous applications we refer the reader to the review articles [15] and [21]. We also refer to the monographs [8] and [9] as the main sources for general facts concerning the theory of Toeplitz matrices and operators. The asymptotic results concerning the Hankel determinants and their applications – both recent and classical, are featured in the papers [10, 11, 15, 19, 20] and in the references therein.

There has been a growing interest in the asymptotics of Toeplitz+Hankel determinants in recent years. A Toeplitz+Hankel matrix is naturally the sum $T_n[\phi; r] + H_n[w; s]$, and thus it has a determinant of the form

$$D_n(\phi, w; r, s) := \det \begin{pmatrix} \phi_r + w_s & \phi_{r-1} + w_{s+1} & \cdots & \phi_{r-n+1} + w_{s+n-1} \\ \phi_{r+1} + w_{s+1} & \phi_r + w_{s+2} & \cdots & \phi_{r-n+2} + w_{s+n} \\ \vdots & \vdots & \ddots & \vdots \\ \phi_{r+n-1} + w_{s+n-1} & \phi_{r+n-2} + w_{s+n} & \cdots & \phi_r + w_{s+2n-2} \end{pmatrix}, \quad (1.3)$$

$r, s \in \mathbb{Z}$, where, naturally, ϕ and w respectively denote the Toeplitz and Hankel symbols. Although there are no results in the literature for the Toeplitz+Hankel determinants where w is supported on the line, the case where w is supported on the unit circle has been considered under specific assumptions. E. Basor and T. Ehrhardt have studied different aspects of these determinants in a series of papers [2, 3, 4, 5, 6] via operator-theoretic tools over the last 20 years or so. In [13], the Riemann–Hilbert technique which has already been proven very effective to study the asymptotics of Toeplitz and Hankel determinants was extended for the first time to the determinants of Toeplitz+Hankel matrices generated by the same symbol $w = \phi$, where the Hankel weight is supported on \mathbb{T} . In that work the symbol was assumed to be of Fisher–Hartwig type and it was further required that the symbol be *even*, i.e., $w = \tilde{w}$.¹ In [2], by employing the relevant results in [13], the authors managed for the first time to find the asymptotics of Toeplitz+Hankel determinants for certain *non-coinciding* symbols. Indeed, they considered

$$\phi(z) = c(z)\phi_0(z), \quad \text{and} \quad w(z) = c(z)d(z)w_0(z), \quad (1.4)$$

where the functions c and d are assumed to be smooth and nonvanishing on the unit circle with zero winding number. Neither c nor d are assumed to be even functions but it is further required that d satisfies the conditions $d\tilde{d} = 1$ (on the unit circle) and $d(\pm 1) = 1$. Furthermore, ϕ_0 is assumed to be an even function of FH type and w_0 is related to ϕ_0 in one of the following four ways: a) $w_0(z) = \pm\phi_0(z)$, b) $w_0(z) = z\phi_0(z)$ and c) $w_0(z) = -z^{-1}\phi_0(z)$.

Since the Riemann–Hilbert analysis carried out in [13] does not allow for different symbols, the primary goal of this paper is to develop a Riemann–Hilbert framework for asymptotic analysis of Toeplitz+Hankel determinants $D_n(\phi, w; r, s)$ where ϕ and w are not a priori related, at least at the level of formulation of the problem. Indeed, asymptotics of Toeplitz+Hankel determinants with different symbols are interesting for several reasons that prompts this research project. For example, the type (1.4) of Toeplitz+Hankel determinants has appeared in the very recent

¹**Notation.** Throughout the paper we will frequently use the notation $\tilde{f}(z)$, to denote $f(z^{-1})$.

work [12] in connection with the analysis of Ising model on the 45° rotated half-plane, or the so-called *zig-zag* half-plane.

Perhaps, our most important motivation behind studying Toeplitz+Hankel determinants is to study the large n asymptotics of the eigenvalues of the Hankel matrix $H_n[w]$ associated to the symbol w . Specifically, we want to extend the recent results [14] concerning the spectral asymptotics of the Toeplitz matrices² to the Hankel case. The key feature which allows an effective asymptotic spectral analysis of Toeplitz matrices and, in particular, the use of the Riemann–Hilbert method, is that the characteristic polynomial of a Toeplitz matrix is again a Toeplitz determinant with a symbol of general Fisher–Hartwig type (i.e., no conditions on the β -parameters). The asymptotics of such Toeplitz determinants is given by the Basor–Tracy formula (first conjectured by E. Basor and C. Tracy and then proved in [13]). However, in the case of Hankel matrices, and this is the crux of the matter, their characteristic polynomials are not Hankel determinants. Indeed, the characteristic polynomial $\det(H_n[w] - \lambda I)$ of the Hankel matrix $H_n[w]$ is a particular *Toeplitz+Hankel* determinant, with $\phi(z) \equiv -\lambda$. Clearly in the case of characteristic polynomial of a Hankel determinant, there is no relationship between ϕ and w , so to study the asymptotics of this determinant, one can not refer to the works [13] or [2] mentioned above. Here again we are directed to a methodological issue which has to be addressed at a fundamental level by formulation of a suitable Riemann–Hilbert problem.

In this paper, we are proposing a version of the Riemann–Hilbert formalism for the asymptotic analysis of Toeplitz+Hankel determinants based on a certain 4×4 Riemann–Hilbert problem. When the Hankel symbol is supported on the unit circle, we introduce the following system of monic orthogonal polynomials $\{\mathcal{P}_n(z)\}$, $\deg \mathcal{P}_n(z) = n$, associated to $D_n(\phi, w; r, s)$:

$$\int_{\mathbb{T}} \mathcal{P}_n(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \mathcal{P}_n(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} = h_n \delta_{n,k}, \quad k = 0, 1, \dots, n.$$

We also show that for $r = s = 1$, if the symbols are analytic in a neighborhood of the unit circle, one can proceed with a 4×4 analogue of the Deift–Zhou non-linear steepest descent method and arrive at a 4×4 *model Riemann–Hilbert problem* on the unit circle which does not contain the parameter n . It is significant to note that one arrives at *the same* model Riemann–Hilbert in the fundamentally different case where w is supported on the interval $[a, b]$, with $0 < a < b < 1$. In this situation we consider the following system of monic orthogonal polynomials $\{P_n(z)\}$, $\deg P_n(z) = n$, associated to $D_n(\phi, w; r, s)$:

$$\int_{\mathbb{T}} P_n(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_a^b P_n(x) x^{k+s} w(x) dx = h_n \delta_{n,k}, \quad k = 0, 1, \dots, n.$$

In this case, we can proceed with the Riemann–Hilbert analysis with $r = 1$ and an *arbitrary* value for $s \in \mathbb{Z}$.

We have been able to solve the model problem for the class of symbols (1.4) considered in [2], in the absence of Fisher–Hartwig singularities. It is important to discuss the relevancy of the two conditions assumed to be satisfied by the function d in [2], in our Riemann–Hilbert framework (see (1.4) and below). Unlike the condition $d(e^{i\theta})d(e^{-i\theta}) = 1$ which is, remarkably, a simplifying condition for the factorization of the model Riemann–Hilbert problem, it should be noticed that the condition $d(\pm 1) = 1$ is not required in the entirety of our Riemann–Hilbert approach. Solving the model problem allows us to find the asymptotics for the norm h_n of the associated orthogonal polynomials. We provide the details of this calculation for the case where w is supported on the unit circle.

In what follows we formulate our main asymptotic result. To do so we need to introduce some notations. Throughout the paper, we will occasionally refer to a symbol f as a *Szegő-type*

²The large n behavior of the individual eigenvalues of Toeplitz matrices has been also addressed in a number of works – see [7] and references therein.

symbol, if a) it is smooth and nonzero on the unit circle, b) has no winding number, and c) admits an analytic continuation in a neighborhood of the unit circle. Also for a given function f , and an oriented contour Γ , we write $f_+(z)$ (resp. $f_-(z)$) to denote the limiting value of $f(\zeta)$, as ζ approaches $z \in \Gamma$ from the left (resp. right) hand side of the oriented contour Γ with respect to its orientation. We also note that the expressions like $\tilde{f}_\pm(z)$, should be understood as the boundary values of the function \tilde{f} at z and should not be confused with the boundary value of the function f at the point $\frac{1}{z}$, in other words, the operation of $z \mapsto \frac{1}{z}$ precedes that of taking the boundary values. Given the Szegő-type symbols $\phi(z)$ and $w(z) = \tilde{d}(z)\phi(z)$, we define

$$\begin{aligned} \alpha(z) &= \exp \left[\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(\phi(\tau))}{\tau - z} d\tau \right], & \beta(z) &= \exp \left[\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(d(\tau))}{\tau - z} d\tau \right], \\ C_\rho(z) &= -\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{1}{\beta_-(\tau)\beta_+(\tau)\tilde{\alpha}_-(\tau)\alpha_+(\tau)(\tau - z)} d\tau, \\ g_{23}(z) &= -\frac{\alpha(0)\tilde{d}(z)\beta(z)}{\tilde{\alpha}(z)}, & g_{43}(z) &= -\alpha^2(0)\beta(z) \left(\frac{\alpha(z)}{\tilde{\phi}(z)} + \frac{\tilde{d}(z)C_\rho(z)}{\tilde{\alpha}(z)} \right), \\ R_{1,23}(z; n) &= \frac{1}{2\pi i} \int_{\Gamma'_i} \frac{\mu^n g_{23}(\mu)}{\mu - z} d\mu, & R_{1,43}(z; n) &= \frac{1}{2\pi i} \int_{\Gamma'_i} \frac{\mu^n g_{43}(\mu)}{\mu - z} d\mu, \end{aligned} \quad (1.5)$$

and finally

$$\mathcal{E}(n) = \frac{2}{\alpha(0)} R_{1,43}(0; n) - C_\rho(0) R_{1,23}(0; n). \quad (1.6)$$

In (1.5), the contour Γ'_i is a circle, oriented counter-clockwise, with radius $r' < 1$ so that the functions ϕ and d are analytic in the annulus $\{z: r' \leq |z| < 1\}$.

Theorem 1.1. *Suppose that $\phi(e^{i\theta})$ is smooth and nonzero on the unit circle with zero winding number, which admits an analytic continuation in a neighborhood of the unit circle. Let $w = d\phi$, where d satisfies all the properties of ϕ in addition to $d(e^{i\theta})d(e^{-i\theta}) = 1$, for all $\theta \in [0, 2\pi)$. Let also*

$$U_0 := \{z: r_i < |z| < r_o: 0 < r_i < 1 < r_o\}, \quad (1.7)$$

be the neighborhood of the unit circle where both functions, $\phi(z)$ and $d(z)$ are analytic. Denote

$$r_0 := \max \{r_i, r_o^{-1}\}, \quad (1.8)$$

and suppose now that there exists such $C > 0$ that for sufficiently large n ,

$$|\mathcal{E}(n)| \geq Cr^n \quad (1.9)$$

for some $r \in [r_0, 1)$, where $\mathcal{E}(n)$ is the functional of the weights ϕ and w defined in (1.6). Then, for sufficiently large n the determinant $D_n(\phi, w; 1, 1) \neq 0$ and the asymptotics of

$$h_{n-1} \equiv \frac{D_n(\phi, w; 1, 1)}{D_{n-1}(\phi, w; 1, 1)},$$

is given by

$$h_{n-1} = -\alpha(0) \frac{\mathcal{E}(n)}{\mathcal{E}(n-1)} (1 + \mathcal{O}(e^{-c_1 n})), \quad n \rightarrow \infty, \quad (1.10)$$

where $c_1 = -\log\left(\frac{r_1^2}{r}\right) > 0$, and r_1 is any number satisfying the conditions: $r < r_1 < 1$ and $r_1^2 < r$.

Remark 1.2. As a slight generalization of Theorem 1.1, one could replace the condition (1.9) by

$$|\mathcal{E}(n)| \geq Cr^n \varkappa(n) \quad (1.11)$$

for some $r \in [r_0, 1)$, where $\varkappa(n) > 0$ for all sufficiently large n and is such that $e^{-\varepsilon n} = o(\varkappa(n))$ for all $\varepsilon > 0$ as $n \rightarrow \infty$. So as a result, (1.10) should be replaced by

$$h_{n-1} = -\alpha(0) \frac{\mathcal{E}(n)}{\mathcal{E}(n-1)} \left(1 + \mathcal{O}\left(\frac{e^{-c_1 n}}{\varkappa(n)}\right) \right), \quad n \rightarrow \infty,$$

where the constant c_1 above is the same as the one in (1.10).

Remark 1.3. The analyticity of ϕ in the neighborhood of the unit circle is a technical condition. It can be lifted and replaced by certain smoothness conditions using the approximation type arguments similar to the ones used in [13, Section 6.2]. However, because of the condition (1.9), the corresponding analysis becomes more subtle. We shall address this issue together with several other technical points in the forthcoming publications.

1.1 Outline

In Section 2 we will analyze the case where Hankel symbol w is supported on the unit circle. We will propose a 2×2 Riemann–Hilbert problem *with a shift* for the associated orthogonal polynomials. For an effective Riemann–Hilbert analysis, we will then propose a 2×4 Riemann–Hilbert problem whose jump conditions could be written in the usual form of matrix multiplications. We will then formulate a 4×4 Riemann–Hilbert problem which is the suitable framework for our analysis. In the formulation of the 4×4 Riemann–Hilbert problem, for technical reasons, we will restrict to the particular offset values $r = s = 1$. Following the natural steps of steepest descent analysis we will arrive at a model problem in Section 2.6, which we will refer to as the *model Riemann–Hilbert problem for the pair* (ϕ, w) .

Section 3 is devoted to analysis of the case where the Hankel symbol is supported on the interval $[a, b]$, with $0 < a < b < 1$. Similar to Section 2, we will propose a 2×2 Riemann–Hilbert problem with a shift for the associated system of orthogonal polynomials and for the same reasons mentioned above we pass through a 2×4 to arrive at a suitable 4×4 Riemann–Hilbert problem. In this case, our methods allow for considering an arbitrary offset for the Hankel part, more precisely, we can pursue the steepest descent analysis for $r = 1$ and an arbitrary $s \in \mathbb{Z}$. This steepest descent analysis leads us to a model problem, which is the same as the model problem of Section 2, except that it is for the pair $(\phi, -\tilde{u})$, where

$$u(z) = z \int_a^b \frac{x^{s-1} w(x)}{x - z} dx.$$

Although there are similarities between the steepest descent analysis of Section 2 and Section 3, we feel obliged to lay out a thorough exposition to illustrate the remarkable fact that the same model RH problem emerges in both cases.

In Section 4, we present the factorization of the model RH problem for the pair $(\phi, d\phi)$, where the functions ϕ and d are of Szegő-type. The function d further satisfies $d\tilde{d} = 1$ on the unit circle. We will then use this solution to construct the solutions to the global parametrix and the small-norm Riemann–Hilbert problems, which finally enables us to find the asymptotics of the norms h_n of the associated monic orthogonal polynomials. Although we do not formulate our main asymptotic result for the symbol pairs discussed in Section 3, in Remark 4.2 we justify that the analysis of Section 4 is completely relevant for the symbol pairs of Section 3 as well.

Finally, in Section 5 we summarize the still open technical and conceptual questions which we are going to address in our future work.

2 Toeplitz+Hankel determinants: Hankel weight supported on \mathbb{T}

In this section we assume that w is supported on the unit circle and that both symbols ϕ and w , admit analytic continuations to a neighborhood of the unit circle. A key observation is that the determinant (1.3) is related to the system of monic polynomials $\{\mathcal{P}_n(z)\}$, $\deg \mathcal{P}_n(z) = n$, determined by the orthogonality relations

$$\int_{\mathbb{T}} \mathcal{P}_n(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \mathcal{P}_n(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} = h_n \delta_{n,k}, \quad k = 0, 1, \dots, n. \quad (2.1)$$

These polynomials exist and are unique if the Toeplitz+Hankel determinants (1.3) are non-zero. This can be seen as follows. Expectedly, if $D_n \equiv D_n(\phi, w; r, s) \neq 0$, the polynomials \mathcal{P}_n can be written as the following determinants

$$\mathcal{P}_n(z) := \frac{1}{D_n} \det \begin{pmatrix} \phi_r + w_s & \phi_{r-1} + w_{s+1} & \cdots & \phi_{r-n} + w_{s+n} \\ \phi_{r+1} + w_{s+1} & \phi_r + w_{s+2} & \cdots & \phi_{r-n+1} + w_{s+n+1} \\ \vdots & \vdots & \ddots & \vdots \\ \phi_{r+n-1} + w_{s+n-1} & \phi_{r+n-2} + w_{s+n} & \cdots & \phi_{r-1} + w_{s+2n-1} \\ 1 & z & \cdots & z^n \end{pmatrix}. \quad (2.2)$$

Indeed, for the polynomials defined by (2.2) we have that

$$\begin{aligned} & \int_{\mathbb{T}} \mathcal{P}_n(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} \\ &= \frac{1}{D_n} \det \begin{pmatrix} \phi_r + w_s & \phi_{r-1} + w_{s+1} & \cdots & \phi_{r-n} + w_{s+n} \\ \phi_{r+1} + w_{s+1} & \phi_r + w_{s+2} & \cdots & \phi_{r-n+1} + w_{s+n+1} \\ \vdots & \vdots & \ddots & \vdots \\ \phi_{r+n-1} + w_{s+n-1} & \phi_{r+n-2} + w_{s+n} & \cdots & \phi_{r-1} + w_{s+2n-1} \\ w_{k+s} & w_{k+s+1} & \cdots & w_{k+s+n} \end{pmatrix}, \end{aligned}$$

and

$$\begin{aligned} & \int_{\mathbb{T}} \mathcal{P}_n(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} \\ &= \frac{1}{D_n} \det \begin{pmatrix} \phi_r + w_s & \phi_{r-1} + w_{s+1} & \cdots & \phi_{r-n} + w_{s+n} \\ \phi_{r+1} + w_{s+1} & \phi_r + w_{s+2} & \cdots & \phi_{r-n+1} + w_{s+n+1} \\ \vdots & \vdots & \ddots & \vdots \\ \phi_{r+n-1} + w_{s+n-1} & \phi_{r+n-2} + w_{s+n} & \cdots & \phi_{r-1} + w_{s+2n-1} \\ \phi_{k+r} & \phi_{k+r-1} & \cdots & \phi_{k+r-n} \end{pmatrix}, \end{aligned}$$

hence

$$\begin{aligned} & \int_{\mathbb{T}} \mathcal{P}_n(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \mathcal{P}_n(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} \\ &= \frac{1}{D_n} \det \begin{pmatrix} \phi_r + w_s & \phi_{r-1} + w_{s+1} & \cdots & \phi_{r-n} + w_{s+n} \\ \phi_{r+1} + w_{s+1} & \phi_r + w_{s+2} & \cdots & \phi_{r-n+1} + w_{s+n+1} \\ \vdots & \vdots & \ddots & \vdots \\ \phi_{r+n-1} + w_{s+n-1} & \phi_{r+n-2} + w_{s+n} & \cdots & \phi_{r-1} + w_{s+2n-1} \\ \phi_{k+r} + w_{k+s} & \phi_{k+r-1} + w_{k+s+1} & \cdots & \phi_{k+r-n} + w_{k+s+n} \end{pmatrix} = \frac{D_{n+1}}{D_n} \delta_{n,k}. \end{aligned}$$

The uniqueness of the polynomial $\mathcal{P}_n(z) = z^n + a_{n-1}z^{n-1} + \dots + a_0$ satisfying (2.1), simply follows from the fact that one has the following linear system for the coefficients $a_j, 0 \leq j \leq n-1$:

$$(T_n[\phi; r] + H_n[w; s]) \begin{pmatrix} a_0 \\ a_1 \\ \vdots \\ a_{n-1} \end{pmatrix} = \begin{pmatrix} -\phi_{-n+r} - w_{n+s} \\ -\phi_{1-n+r} - w_{1+n+s} \\ \vdots \\ -\phi_{-1+r} - w_{2n-1+s} \end{pmatrix}. \quad (2.3)$$

So if $D_n \neq 0$, the coefficients a_j and hence \mathcal{P}_n , can be uniquely determined by inverting the Toeplitz+Hankel matrix in (2.3). So the polynomials defined by (2.2) are the unique polynomials satisfying (2.1), and

$$h_n = \frac{D_{n+1}(\phi, w; r, s)}{D_n(\phi, w; r, s)}. \quad (2.4)$$

We consider the function \mathcal{Y} defined as

$$\mathcal{Y}(z; n) = \begin{pmatrix} \mathcal{P}_n(z) & \int_{\mathbb{T}} \frac{\xi^s \tilde{w}(\xi) \mathcal{P}_n(\xi) + \xi^r \tilde{\phi}(\xi) \tilde{\mathcal{P}}_n(\xi)}{\xi - z} \frac{d\xi}{2\pi i \xi} \\ -\frac{1}{h_{n-1}} \mathcal{P}_{n-1}(z) & -\frac{1}{h_{n-1}} \int_{\mathbb{T}} \frac{\xi^s \tilde{w}(\xi) \mathcal{P}_{n-1}(\xi) + \xi^r \tilde{\phi}(\xi) \tilde{\mathcal{P}}_{n-1}(\xi)}{\xi - z} \frac{d\xi}{2\pi i \xi} \end{pmatrix}, \quad (2.5)$$

assuming that $D_n, D_{n-1} \neq 0$, which ensures that a) $h_{n-1} = D_n/D_{n-1}$ is well-defined and nonzero and b) that \mathcal{P}_n and \mathcal{P}_{n-1} exist and are unique. Now consider the following Riemann–Hilbert problem for finding the 2×2 matrix \mathcal{Y} satisfying:

- **RH- $\mathcal{Y}1$** \mathcal{Y} is holomorphic in $\mathbb{C} \setminus \mathbb{T}$.
- **RH- $\mathcal{Y}2$** For $z \in \mathbb{T}$ we have

$$\mathcal{Y}_+^{(1)}(z; n) = \mathcal{Y}_-^{(1)}(z; n), \quad z \in \mathbb{T}, \quad (2.6)$$

and

$$\mathcal{Y}_+^{(2)}(z; n) = \mathcal{Y}_-^{(2)}(z; n) + z^{-1+s} \tilde{w}(z) \mathcal{Y}_-^{(1)}(z; n) + z^{-1+r} \tilde{\phi}(z) \mathcal{Y}_-^{(1)}(z^{-1}; n), \quad z \in \mathbb{T}, \quad (2.7)$$

where \mathbb{T} is positively oriented in the counter-clockwise direction.

- **RH- $\mathcal{Y}3$** As $z \rightarrow \infty$, \mathcal{Y} satisfies

$$\mathcal{Y}(z; n) = (I + \mathcal{O}(z^{-1})) z^{n\sigma_3} = \begin{pmatrix} z^n + \mathcal{O}(z^{n-1}) & \mathcal{O}(z^{n-1}) \\ \mathcal{O}(z^{n-1}) & z^{-n} + \mathcal{O}(z^{-n-1}) \end{pmatrix}, \quad (2.8)$$

where $\mathcal{Y}^{(1)}$ and $\mathcal{Y}^{(2)}$ are the first and second columns of \mathcal{Y} respectively and

$$\sigma_3 = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}$$

is the third Pauli matrix. Notice that the jump conditions for the \mathcal{Y} -RHP can not be written in the usual form of matrix multiplications. Expectedly, this feature of the \mathcal{Y} -RHP (which will be addressed in more detail in Section 2.1) significantly affects the progression of the Riemann–Hilbert analysis as it necessitates studying larger size Riemann–Hilbert problems.

The next theorem establishes the association of the function (2.5) with the \mathcal{Y} -RHP.

Theorem 2.1. *The following statements are true.*

1. Suppose that $D_n, D_{n-1} \neq 0$. Then, the Riemann–Hilbert problem **RH- $\mathcal{Y}1$** through **RH- $\mathcal{Y}3$** is uniquely solvable and its solution \mathcal{Y} is defined by (2.5). Moreover,

$$h_{n-1} = - \lim_{z \rightarrow \infty} z^{n-1} / \mathcal{Y}_{21}(z; n). \quad (2.9)$$

2. Suppose that the Riemann–Hilbert problem **RH- $\mathcal{Y}1$** through **RH- $\mathcal{Y}3$** has a unique solution. Then $D_n \neq 0$, $\text{rank}(T_{n-1}[\phi; r] + H_{n-1}[w; s]) \geq n - 2$, and $\mathcal{P}_n(z) = \mathcal{Y}_{11}(z; n)$.
3. Suppose that the Riemann–Hilbert problem **RH- $\mathcal{Y}1$** through **RH- $\mathcal{Y}3$** has a unique solution. Suppose also that

$$\lim_{z \rightarrow \infty} \mathcal{Y}_{21}(z; n) z^{-n+1} \neq 0.$$

Then, as before, $D_n \neq 0$, $\mathcal{P}_n(z) = \mathcal{Y}_{11}(z; n)$, and, in addition,

$$D_{n-1} \neq 0, \quad h_{n-1} = - \lim_{z \rightarrow \infty} \mathcal{Y}_{21}^{-1}(z; n) z^{n-1}, \quad \mathcal{P}_{n-1}(z) = -h_{n-1} \mathcal{Y}_{21}(z; n).$$

Proof. Assume that $D_n, D_{n-1} \neq 0$, and thus the function \mathcal{Y} is uniquely defined by (2.5), being identified with the *unique* orthogonal polynomials satisfying the orthogonality conditions (2.1). We now show that \mathcal{Y} given by (2.5) is a unique solution of **RH- $\mathcal{Y}1$** through **RH- $\mathcal{Y}3$** . It is clear that **RH- $\mathcal{Y}1$** is satisfied due to general properties of Cauchy integrals. From (2.6) we see that \mathcal{Y}_{11} and \mathcal{Y}_{21} are entire functions, and from (2.8) we know that \mathcal{Y}_{11} has to be a monic polynomial of degree n and \mathcal{Y}_{21} has to be a polynomial of degree $n - 1$ or less,

$$\mathcal{Y}_{11}(z; n) = z^n + \sum_{k=0}^{n-1} \alpha_k z^k, \quad \mathcal{Y}_{21}(z; n) = \sum_{k=0}^{n-1} \beta_k z^k.$$

From (2.7) and what we just mentioned about \mathcal{Y}_{11} we would have

$$\mathcal{Y}_{12,+}(z; n) - \mathcal{Y}_{12,-}(z; n) = z^{-1+s} \tilde{w}(z) \mathcal{Y}_{11}(z; n) + z^{-1+r} \tilde{\phi}(z) \tilde{\mathcal{Y}}_{11}(z; n).$$

So by Plemelj–Sokhotskii formula we have

$$\mathcal{Y}_{12}(z; n) = \frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\xi^{-1+s} \tilde{w}(\xi) \mathcal{Y}_{11}(\xi; n) + \xi^{-1+r} \tilde{\phi}(\xi) \tilde{\mathcal{Y}}_{11}(\xi; n)}{\xi - z} d\xi.$$

Using the identity

$$\frac{1}{\xi - z} = - \sum_{k=0}^n \frac{\xi^k}{z^{k+1}} + \frac{\xi^{n+1}}{(\xi - z) z^{n+1}},$$

we get

$$\begin{aligned} \mathcal{Y}_{12}(z; n) &= - \sum_{k=0}^n \frac{1}{z^{k+1}} \int_{\mathbb{T}} [\xi^{-1+s} \tilde{w}(\xi) \mathcal{Y}_{11}(\xi; n) \xi^k + \xi^{-1+r} \tilde{\phi}(\xi) \tilde{\mathcal{Y}}_{11}(\xi; n) \xi^k] \frac{d\xi}{2\pi i} \\ &\quad + \frac{1}{z^{n+1}} \int_{\mathbb{T}} \frac{\xi^{n+1}}{(\xi - z)} [\xi^{-1+s} \tilde{w}(\xi) \mathcal{Y}_{11}(\xi; n) + \xi^{-1+r} \tilde{\phi}(\xi) \tilde{\mathcal{Y}}_{11}(\xi; n)] \frac{d\xi}{2\pi i}. \end{aligned}$$

Note that since $\mathcal{Y}_{12}(z; n) = O(z^{-n-1})$, we must have:

$$\int_{\mathbb{T}} \tilde{w}(\xi) \mathcal{Y}_{11}(\xi; n) \xi^{k+s} \frac{d\xi}{2\pi i \xi} + \int_{\mathbb{T}} \tilde{\phi}(\xi) \tilde{\mathcal{Y}}_{11}(\xi; n) \xi^{k+r} \frac{d\xi}{2\pi i \xi} = 0, \quad 0 \leq k \leq n - 1.$$

In the second integral we make the change of variable $\xi \mapsto \tau := \xi^{-1}$ and as a result we will arrive at

$$\int_{\mathbb{T}} \mathcal{Y}_{11}(\xi; n) \xi^{k+s} \tilde{w}(\xi) \frac{d\xi}{2\pi i \xi} + \int_{\mathbb{T}} \mathcal{Y}_{11}(\tau; n) \tau^{-k-r} \phi(\tau) \frac{d\tau}{2\pi i \tau} = 0, \quad 0 \leq k \leq n-1.$$

Since \mathcal{Y}_{11} satisfies the orthogonality relations (2.1) we necessarily have

$$\mathcal{Y}_{11}(z; n) = \mathcal{P}_n(z).$$

In a similar fashion one can show that

$$\mathcal{Y}_{22}(z; n) = \frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\xi^{-1+s} \tilde{w}(\xi) \mathcal{Y}_{21}(\xi; n) + \xi^{-1+r} \tilde{\phi}(\xi) \tilde{\mathcal{Y}}_{21}(\xi; n)}{\xi - z} d\xi.$$

The asymptotic condition, $\mathcal{Y}_{22}(z; n) = z^{-n} + O(z^{-n-1})$, would then yield the orthogonality relations,

$$\int_{\mathbb{T}} \mathcal{Y}_{21}(\xi; n) \xi^{k+s} \tilde{w}(\xi) \frac{d\xi}{2\pi i \xi} + \int_{\mathbb{T}} \mathcal{Y}_{21}(\tau; n) \tau^{-k-r} \phi(\tau) \frac{d\tau}{2\pi i \tau} = -\delta_{k, n-1}, \quad 0 \leq k \leq n-1.$$

These relations in turn are equivalent to the following linear system for the coefficients β_j ,

$$(T_n[\phi; r] + H_n[w; s]) \begin{pmatrix} \beta_0 \\ \beta_1 \\ \vdots \\ \beta_{n-1} \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ \vdots \\ -1 \end{pmatrix}.$$

From this it follows that β_{n-1} is necessarily not zero. Indeed, using Cramer's rule, we arrive at the formula,

$$\beta_{n-1} = -\frac{D_{n-1}}{D_n} \neq 0.$$

This would also mean that $\beta_{n-1} = -\frac{1}{h_{n-1}}$ and

$$\mathcal{Y}_{21}(z; n) = -\frac{1}{h_{n-1}} \mathcal{P}_{n-1}(z). \quad (2.10)$$

This finishes the argument why (2.5) uniquely satisfies both **RH- $\mathcal{Y}2$** and **RH- $\mathcal{Y}3$** as well. Also, equation (2.10) implies equation (2.9).

To prove the second statement, now assume that the Riemann–Hilbert problem **RH- $\mathcal{Y}1$** through **RH- $\mathcal{Y}3$** has a unique solution. Because of similar considerations mentioned above, it can be written as

$$\mathcal{Y}(z; n) = \begin{pmatrix} \mathcal{Q}_n(z) & \int_{\mathbb{T}} \frac{\xi^s \tilde{w}(\xi) \mathcal{Q}_n(\xi) + \xi^r \tilde{\phi}(\xi) \tilde{\mathcal{Q}}_n(\xi)}{\xi - z} \frac{d\xi}{2\pi i \xi} \\ \mathcal{R}_{n-1}(z) & \int_{\mathbb{T}} \frac{\xi^s \tilde{w}(\xi) \mathcal{R}_{n-1}(\xi) + \xi^r \tilde{\phi}(\xi) \tilde{\mathcal{R}}_{n-1}(\xi)}{\xi - z} \frac{d\xi}{2\pi i \xi} \end{pmatrix}, \quad (2.11)$$

where $\mathcal{Q}_n(z)$ and $\mathcal{R}_m(z)$ are polynomials,

$$\mathcal{Q}_n(z) = z^n + \sum_{k=0}^{n-1} \alpha_k z^k, \quad \mathcal{R}_{n-1}(z) = \sum_{k=0}^{n-1} \beta_k z^k,$$

satisfying the orthogonality conditions,

$$\int_{\mathbb{T}} \mathcal{Q}_n(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \mathcal{Q}_n(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} = 0, \quad k = 0, 1, \dots, n-1, \quad (2.12)$$

and

$$\int_{\mathbb{T}} \mathcal{R}_{n-1}(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \mathcal{R}_{n-1}(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} = -\delta_{k,n-1}, \quad (2.13)$$

$$k = 0, 1, \dots, n-1.$$

Note that the above equations can be written as the following linear systems on the coefficients α_j and β_j :

$$(T_n[\phi; r] + H_n[w; s]) \begin{pmatrix} \alpha_0 \\ \alpha_1 \\ \vdots \\ \alpha_{n-1} \end{pmatrix} = \begin{pmatrix} -\phi_{-n+r} - w_{n+s} \\ -\phi_{1-n+r} - w_{1+n+s} \\ \vdots \\ -\phi_{-1+r} - w_{2n-1+s} \end{pmatrix}, \quad (2.14)$$

and

$$(T_n[\phi; r] + H_n[w; s]) \begin{pmatrix} \beta_0 \\ \beta_1 \\ \vdots \\ \beta_{n-1} \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ \vdots \\ -1 \end{pmatrix}. \quad (2.15)$$

Note that since we have assumed that the solution of \mathcal{Y} -RHP exists, both linear systems (2.14) and (2.15) must have solutions; and since the solution of \mathcal{Y} -RHP is unique, it follows that the determinant of the systems (2.14) and (2.15), which is nothing but D_n , is nonzero. Thus, we certainly have that

$$\mathcal{Q}_n(z) = \mathcal{P}_n(z).$$

It is now obvious that $\text{rank}(T_{n-1}[\phi; r] + H_{n-1}[w; s]) \geq n-2$. Because if $D_{n-1} = 0$ and $D_n \neq 0$, one can easily show that $\text{rank}(T_{n-1}[\phi; r] + H_{n-1}[w; s]) = n-2$. Although the proof of the second statement is complete, it is worthwhile to say more about \mathcal{Y}_{21} when $D_{n-1} = 0$, because then the objects $\mathcal{P}_{n-1}(z)$ and h_{n-1} do not exist. Indeed the homogenous system,

$$(T_{n-1}[\phi; r] + H_{n-1}[w; s]) \begin{pmatrix} c_0 \\ c_1 \\ \vdots \\ c_{n-2} \end{pmatrix} = 0, \quad (2.16)$$

has infinitely many nontrivial solutions. Let us take one of them and define the following polynomial of degree $n-2$

$$\hat{\mathcal{R}}_{n-2}(z) = \sum_{j=0}^{n-2} c_j z^j,$$

which automatically satisfies

$$\int_{\mathbb{T}} \hat{\mathcal{R}}_{n-2}(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \hat{\mathcal{R}}_{n-2}(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} = 0, \quad k = 0, 1, \dots, n-2.$$

Put

$$\delta_{n-1} := \int_{\mathbb{T}} \hat{\mathcal{R}}_{n-2}(z) z^{-n+1-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \hat{\mathcal{R}}_{n-2}(z) z^{n-1+s} \tilde{w}(z) \frac{dz}{2\pi i z}.$$

Note that $\delta_{n-1} \neq 0$. Otherwise, we could redefine $\mathcal{R}_{n-1}(z)$ as

$$\mathcal{R}_{n-1}(z) \rightarrow \mathcal{R}_{n-1}(z) + \hat{\mathcal{R}}_{n-2}(z),$$

which, on the contrary, means that the solution of the \mathcal{Y} -RHP is not unique. Hence, we can find a *unique* normalization of the nontrivial solution of (2.16), to assure that

$$\delta_{n-1} = -1.$$

This, together with the uniqueness of the solution of the \mathcal{Y} -RHP implies that $\mathcal{R}_{n-1}(z) = \hat{\mathcal{R}}_{n-2}(z)$.

To prove the third statement, we notice that, due to the additional condition at $z = \infty$, the unique solution of the Riemann–Hilbert problem **RH- \mathcal{Y} 1** through **RH- \mathcal{Y} 3** can be written in the same form (2.11) with polynomial $\mathcal{R}_{n-1}(z)$ replaced by

$$\mathcal{R}_{n-1}(z) = q_n \mathcal{Q}_{n-1}(z), \quad q \neq 0,$$

where both $\mathcal{Q}_n(z)$ and $\mathcal{Q}_{n-1}(z)$ are monic polynomials,

$$\mathcal{Q}_n(z) = z^n + \sum_{k=0}^{n-1} \alpha_k z^k, \quad \mathcal{Q}_{n-1}(z) = z^{n-1} + \sum_{k=0}^{n-2} \beta_k z^k,$$

satisfying the orthogonality conditions,

$$\int_{\mathbb{T}} \mathcal{Q}_n(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \mathcal{Q}_n(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} = 0, \quad k = 0, 1, \dots, n-1,$$

and

$$q_n \left(\int_{\mathbb{T}} \mathcal{Q}_{n-1}(z) z^{-k-r} \phi(z) \frac{dz}{2\pi i z} + \int_{\mathbb{T}} \mathcal{Q}_{n-1}(z) z^{k+s} \tilde{w}(z) \frac{dz}{2\pi i z} \right) = -\delta_{k,n-1}, \quad (2.17)$$

$$k = 0, 1, \dots, n-1.$$

As before with equations (2.12) and (2.13), the above equations can be written as the following linear systems on the coefficients α_j and β_j :

$$(T_n[\phi; r] + H_n[w; s]) \begin{pmatrix} \alpha_0 \\ \alpha_1 \\ \vdots \\ \alpha_{n-1} \end{pmatrix} = \begin{pmatrix} -\phi_{-n+r} - w_{n+s} \\ -\phi_{1-n+r} - w_{1+n+s} \\ \vdots \\ -\phi_{-1+r} - w_{2n-1+s} \end{pmatrix},$$

and

$$(T_{n-1}[\phi; r] + H_{n-1}[w; s]) \begin{pmatrix} \beta_0 \\ \beta_1 \\ \vdots \\ \beta_{n-2} \end{pmatrix} = \begin{pmatrix} -\phi_{-n+1+r} - w_{n-1+s} \\ -\phi_{1-n+1+r} - w_{1+n-1+s} \\ \vdots \\ -\phi_{-1+r} - w_{2n-3+s} \end{pmatrix}.$$

We now argue, that both D_n and D_{n-1} are nonzero and thus $\mathcal{Q}_n = \mathcal{P}_n$ and $\mathcal{Q}_{n-1} = \mathcal{P}_{n-1}$ are given by (2.2). Otherwise, we shall have more than one solution for one or both of the above linear systems, which means that we could find distinct polynomials $\mathcal{Q}_n(z)$ (if $D_n = 0$), or distinct polynomials $\mathcal{Q}_{n-1}(z)$ (if $D_{n-1} = 0$), or both (if $D_n, D_{n-1} = 0$). Either way, using these distinct polynomials, we could construct distinct solutions of the \mathcal{Y} -RHP which contradicts our assumption. This finishes the proof of the theorem. Note that from (2.17) we necessarily have $q_n \neq 0$, and since $\mathcal{Q}_{n-1} = \mathcal{P}_{n-1}$ we conclude that $h_{n-1} = -1/q_n \neq 0$, which could also be seen from $h_{n-1} = D_n/D_{n-1}$. \blacksquare

Corollary 2.2. *Suppose that the \mathcal{Y} -RH problem has a unique solution for n and $n-1$. Then*

$$D_n \neq 0, \quad D_{n-1} \neq 0, \quad \text{and} \quad h_{n-1} \neq 0,$$

where h_{n-1} can be reconstructed from the RHP data as

$$h_{n-1} = - \lim_{z \rightarrow \infty} z^{n-1} / \mathcal{Y}_{21}(z; n). \quad (2.18)$$

Equation (2.18), as in the pure Toeplitz or Hankel cases, in conjunction with (2.4) reduces the asymptotic analysis of the Toeplitz+Hankel determinants to the asymptotic analysis of the Riemann–Hilbert problem for \mathcal{Y} .

2.1 The associated 2×4 and 4×4 Riemann–Hilbert problems

In the rest of this section we will develop a 4×4 analogue of the Deift/Zhou non-linear steepest descent method for the Toeplitz+Hankel determinants (1.3). The \mathcal{Y} -RHP is a particular case of the matrix *Riemann–Hilbert problem with a shift*, or the matrix analytical boundary problem of the *Carleman* type. Indeed, the matrix form of the equations (2.6)–(2.7) reads as follows

$$\mathcal{Y}_+(z; n) = \mathcal{Y}_-(z; n)G_1(z) + \mathcal{Y}_-(\kappa(z); n)G_2(z),$$

where

$$G_1(z) = \begin{pmatrix} 1 & z^{s-1}\tilde{w}(z) \\ 0 & 1 \end{pmatrix}, \quad G_2(z) = \begin{pmatrix} 0 & z^{r-1}\tilde{\phi}(z) \\ 0 & 0 \end{pmatrix},$$

and the “shift” κ is the mapping

$$\kappa(z) = \frac{1}{z}.$$

The presence of the shift makes it impossible to directly apply the usual 2×2 version of the Deift–Zhou nonlinear steepest descent method to the \mathcal{Y} -RHP. However, the mapping κ satisfies the Carleman condition, $\kappa(\kappa(z)) = z$, and hence we can translate the 2×2 \mathcal{Y} -RHP to the usual matrix form by doubling the relevant matrix sizes. More precisely, we first propose the associated 2×4 and then the associated 4×4 Riemann–Hilbert problems. Although more complicated, the analysis of the proposed 4×4 Riemann–Hilbert problem follows in the same spirit as the lower dimensional RHPs until we get to the model Riemann–Hilbert problem for Toeplitz+Hankel determinants introduced in Section 2.6.

Let us define the 2×4 matrix $\overset{\circ}{\mathcal{X}}$ out of the columns of \mathcal{Y} as follows

$$\overset{\circ}{\mathcal{X}}(z; n) := (\mathcal{Y}^{(1)}(z; n), \tilde{\mathcal{Y}}^{(1)}(z; n), \mathcal{Y}^{(2)}(z; n), \tilde{\mathcal{Y}}^{(2)}(z; n)). \quad (2.19)$$

From (2.6), (2.7) and (2.8) we obtain the following Riemann–Hilbert problem for $\overset{\circ}{\mathcal{X}}$:

- **RH- $\overset{\circ}{\mathcal{X}}1$** $\overset{\circ}{\mathcal{X}}$ is holomorphic in $\mathbb{C} \setminus (\mathbb{T} \cup \{0\})$.
- **RH- $\overset{\circ}{\mathcal{X}}2$** For $z \in \mathbb{T}$, $\overset{\circ}{\mathcal{X}}$ satisfies

$$\overset{\circ}{\mathcal{X}}_+(z; n) = \overset{\circ}{\mathcal{X}}_-(z; n) \begin{pmatrix} 1 & 0 & z^{s-1}\tilde{w}(z) & -z^{-r+1}\phi(z) \\ 0 & 1 & z^{r-1}\tilde{\phi}(z) & -z^{-s+1}w(z) \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}. \quad (2.20)$$

- **RH- $\overset{\circ}{\mathcal{X}}3$** As $z \rightarrow \infty$ we have

$$\begin{aligned} \overset{\circ}{\mathcal{X}}(z; n) &= \begin{pmatrix} 1 + \mathcal{O}(z^{-1}) & C_1(n) + \mathcal{O}(z^{-1}) & \mathcal{O}(z^{-1}) & C_3(n) + \mathcal{O}(z^{-1}) \\ \mathcal{O}(z^{-1}) & C_2(n) + \mathcal{O}(z^{-1}) & 1 + \mathcal{O}(z^{-1}) & C_4(n) + \mathcal{O}(z^{-1}) \end{pmatrix} \\ &\quad \times \begin{pmatrix} z^n & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^{-n} & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}. \end{aligned} \quad (2.21)$$

- **RH- $\overset{\circ}{\mathcal{X}}4$** As $z \rightarrow 0$ we have

$$\begin{aligned} \overset{\circ}{\mathcal{X}}(z; n) &= \begin{pmatrix} C_1(n) + \mathcal{O}(z) & 1 + \mathcal{O}(z) & C_3(n) + \mathcal{O}(z) & \mathcal{O}(z) \\ C_2(n) + \mathcal{O}(z) & \mathcal{O}(z) & C_4(n) + \mathcal{O}(z) & 1 + \mathcal{O}(z) \end{pmatrix} \\ &\quad \times \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^{-n} & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^n \end{pmatrix}, \end{aligned} \quad (2.22)$$

where

$$C_1(n) = \mathcal{Y}_{11}(0; n), \quad C_3(n) = \mathcal{Y}_{12}(0; n), \quad C_2(n) = \mathcal{Y}_{21}(0; n), \quad C_4(n) = \mathcal{Y}_{22}(0; n).$$

Remark 2.3. It is well-known that for construction of the global parametrix in the Riemann–Hilbert analysis, one has to construct the so-called Szegő functions of the following type:

$$\mathcal{S}_f(z) := \exp \left[\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(f(\tau))}{\tau - z} d\tau \right],$$

where one assumes that f is non-zero and continuous on the unit circle with zero winding number. The natural progression of the Riemann–Hilbert analysis with general offset values $r, s \in \mathbb{Z}$, leads us to constructing Szegő functions for $f_1(z) = z^{1-r}\phi(z)$ and $f_2(z) = z^{1-s}w(z)$ (or, equivalently, for f_1 and f_2/f_1). Although there are ways to “peel off” the winding generating parts of f_1 and f_2 (see Section 5.2). The first natural case to consider is when f_1 and f_2 have zero winding numbers. To that end, in this work we shall only consider the case $r = s = 1$ and the symbols ϕ and w which are of Szegő type and analytic in a neighborhood of the unit circle. We also discuss prospects of extension of our method to symbols with Fisher–Hartwig singularities in Section 5.4. Also it is worth mentioning that our method in the present work trivially extends to the case where we consider general, but fixed, $r, s \in \mathbb{Z}$ with symbols ϕ and w respectively with winding numbers $r - 1$ and $s - 1$ (so that the overall winding number of f_1 and f_2 are zero, which is the main condition to be considered). It should be also mentioned that in the pure Toeplitz case

there is another, and historically the first, way to put the analysis of Toeplitz determinants in the framework of the Riemann–Hilbert scheme. This approach was suggested by Jinho Baik, Percy Deift, and Kurt Johansson in [1]. The Riemann–Hilbert problem that appear there yields to the Szegő function for $f(z)$ *coinciding* with the original Toeplitz symbol $\phi(z)$. Therefore, the Riemann–Hilbert method of [1] is applied directly to the symbols with zero winding numbers. The crux of the matter is that the construction of [1] can not be extended to the mixed, T+H, situation. The \mathcal{Y} -RHP we are introducing in this paper seems to be the only way to put the mixed problem into the Riemann–Hilbert setting; hence the inevitability of the choice $r = s = 1$, or equivalently, starting with symbols that have nonzero winding numbers.

In view of this remark, in the rest of the paper we assume that $r = s = 1$. In a natural way we now consider the following 4×4 Riemann–Hilbert problem which we will refer to as the \mathcal{X} -RHP.

- **RH- $\mathcal{X}1$** \mathcal{X} is holomorphic in $\mathbb{C} \setminus (\mathbb{T} \cup \{0\})$.
- **RH- $\mathcal{X}2$** For $z \in \mathbb{T}$, \mathcal{X} satisfies

$$\mathcal{X}_+(z; n) = \mathcal{X}_-(z; n) \begin{pmatrix} 1 & 0 & \tilde{w}(z) & -\phi(z) \\ 0 & 1 & \tilde{\phi}(z) & -w(z) \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}. \quad (2.23)$$

- **RH- $\mathcal{X}3$** As $z \rightarrow \infty$ we have

$$\mathcal{X}(z; n) = (I + \mathcal{O}(z^{-1})) \begin{pmatrix} z^n & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^{-n} & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}. \quad (2.24)$$

- **RH- $\mathcal{X}4$** As $z \rightarrow 0$ we have

$$\mathcal{X}(z; n) = P(n)(I + \mathcal{O}(z)) \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^{-n} & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^n \end{pmatrix}. \quad (2.25)$$

Remark 2.4. The uniqueness of the solution of \mathcal{X} -RHP is established using the standard Liouville theorem-based arguments. We also note that the matrix factor $P(n)$ in (2.25) is not a priori prescribed.

Remark 2.5. It is easy to see that the solution \mathcal{X} of the \mathcal{X} -RHP satisfies the symmetry relation,

$$WP^{-1}(n)\mathcal{X}(z^{-1}; n)W = \mathcal{X}(z; n), \quad (2.26)$$

where

$$W = \begin{pmatrix} 0 & 1 & 0 & 0 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 1 & 0 \end{pmatrix}. \quad (2.27)$$

Equation (2.26) in turn yields the following symmetry equation for $P(n)$,

$$P(n) = WP^{-1}(n)W, \quad (2.28)$$

or (taking into account that $W^{-1} = W$),

$$(WP(n))^2 = (P(n)W)^2 = I_4.$$

Remark 2.6. The matrix $P(n)W - I_4$ has rank 2. Here is the proof of this statement which is motivated by some of the referee’s remarks.³

If we take the limit $z \rightarrow 1$, with $|z| > 1$ in (2.26) it follows that

$$WP^{-1}(n)\mathcal{X}_+(1;n)W = \mathcal{X}_-(1;n),$$

or

$$\mathcal{X}_+(1;n)W\mathcal{X}_-^{-1}(1;n) = P(n)W.$$

Using (2.23), the last equation reads,

$$\mathcal{X}_-(1;n)G_{\mathcal{X}}(1)W\mathcal{X}_-^{-1}(1;n) = P(n)W,$$

where $G_{\mathcal{X}}(z)$ is the jump matrix in (2.23). Hence the matrices $P(n)W$ and $G_{\mathcal{X}}(1)W$ are similar, and thus the claim is true if the matrix $G_{\mathcal{X}}(1)W$ has rank 2. We have,

$$G_{\mathcal{X}}(1)W - I_4 = \begin{pmatrix} -1 & 1 & -\phi & w \\ 1 & -1 & -w & \phi \\ 0 & 0 & -1 & 1 \\ 0 & 0 & 1 & -1 \end{pmatrix},$$

where $\phi \equiv \phi(1) = \tilde{\phi}(1)$, and $w \equiv w(1) = \tilde{w}(1)$. The row echelon form of this matrix is

$$\begin{pmatrix} -1 & 1 & -\phi & w \\ 0 & 0 & -1 & 1 \\ 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \end{pmatrix},$$

and the statement about the rank of matrix $G_{\mathcal{X}}(1)W - I_4$ (and hence the rank of matrix $P(n)W - I_4$) follows.

2.2 Relation of the 2×4 and the 4×4 Riemann–Hilbert problems

Put

$$\mathfrak{R}(z;n) := \overset{\circ}{\mathcal{X}}(z;n)\mathcal{X}^{-1}(z;n). \quad (2.29)$$

From (2.20) and (2.23) it is clear that \mathfrak{R} has no jumps. From (2.22) and (2.25) we can obtain the behavior of \mathfrak{R} near zero:

$$\mathfrak{R}(z;n) = \begin{pmatrix} C_1(n) + \mathcal{O}(z) & 1 + \mathcal{O}(z) & C_3(n) + \mathcal{O}(z) & \mathcal{O}(z) \\ C_2(n) + \mathcal{O}(z) & \mathcal{O}(z) & C_4(n) + \mathcal{O}(z) & 1 + \mathcal{O}(z) \end{pmatrix} P^{-1}(n).$$

Therefore \mathfrak{R} is an entire function. Also note that from (2.21) and (2.24) we have

$$\mathfrak{R}(z;n) = \begin{pmatrix} 1 + \mathcal{O}(z^{-1}) & C_1(n) + \mathcal{O}(z^{-1}) & \mathcal{O}(z^{-1}) & C_3(n) + \mathcal{O}(z^{-1}) \\ \mathcal{O}(z^{-1}) & C_2(n) + \mathcal{O}(z^{-1}) & 1 + \mathcal{O}(z^{-1}) & C_4(n) + \mathcal{O}(z^{-1}) \end{pmatrix}, \quad z \rightarrow \infty.$$

Therefore by Liouville’s theorem we conclude that

$$\mathfrak{R}(z;n) = \begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix}. \quad (2.30)$$

³These valuable remarks also led us to formulate and prove Lemmas 2.7 and 2.8, and their analogues in Section 3.

And therefore we have

$$\begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix} = \begin{pmatrix} C_1(n) & 1 & C_3(n) & 0 \\ C_2(n) & 0 & C_4(n) & 1 \end{pmatrix} P^{-1}(n). \quad (2.31)$$

We argue that this system, under certain generic assumptions, is a well-defined linear system on $C_j(n)$ which is uniquely solvable. To see this, let us first double the system, that is consider instead of (2.31), the system

$$\begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix} = \begin{pmatrix} C'_1(n) & 1 & C'_3(n) & 0 \\ C'_2(n) & 0 & C'_4(n) & 1 \end{pmatrix} P^{-1}(n). \quad (2.32)$$

This is an 8×8 system of linear equations for eight unknowns – $C_j(n)$ and $C'_j(n)$. By a straightforward and, in fact, rather simple calculations one finds that the determinant of this 8×8 system is

$$(P_{22}(n)P_{44}(n) - P_{42}(n)P_{24}(n))^2,$$

where $P_{jk}(n)$, $j, k = 1, \dots, 4$ denote the entries of matrix $P(n)$. Hence, assuming that

$$P_{22}(n)P_{44}(n) - P_{42}(n)P_{24}(n) \neq 0, \quad (2.33)$$

we would have the unique solvability of system (2.32). That is, there is only one 8-vector,

$$\vec{C}(n) := (C_1(n), C_2(n), C_3(n), C_4(n), C'_1(n), C'_2(n), C'_3(n), C'_4(n)), \quad (2.34)$$

which solves (2.32). At the same time, the symmetry relation (2.28) would imply that

$$\begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix} W = \begin{pmatrix} C'_1(n) & 1 & C'_3(n) & 0 \\ C'_2(n) & 0 & C'_4(n) & 1 \end{pmatrix} W P(n),$$

or

$$\begin{pmatrix} 1 & C'_1(n) & 0 & C'_3(n) \\ 0 & C'_2(n) & 1 & C'_4(n) \end{pmatrix} = \begin{pmatrix} C_1(n) & 1 & C_3(n) & 0 \\ C_2(n) & 0 & C_4(n) & 1 \end{pmatrix} P^{-1}(n).$$

In other words, together with (2.34), the system (2.32) will be also solved by the vector

$$\vec{C}'(n) := (C'_1(n), C'_2(n), C'_3(n), C'_4(n), C_1(n), C_2(n), C_3(n), C_4(n)).$$

Because of the uniqueness, we conclude that the vectors $\vec{C}'(n)$ and $\vec{C}(n)$ must coincide and hence we must have,

$$C_j(n) = C'_j(n), \quad j = 1, 2, 3, 4.$$

Therefore we have the unique solvability of the original system (2.31) under the generic condition (2.33).

Indeed, we can make the solvability condition of the system (2.31) more flexible. This is shown in the next lemma.

Lemma 2.7. *Assume that at least one of the following six inequalities is true,*

$$P_{22}(n)P_{44}(n) - P_{42}(n)P_{24}(n) \neq 0, \quad (2.35)$$

$$(1 - P_{21}(n))P_{42}(n) + P_{22}(n)P_{41}(n) \neq 0, \quad (2.36)$$

$$(1 - P_{43}(n))P_{22}(n) + P_{23}(n)P_{42}(n) \neq 0, \quad (2.37)$$

$$(1 - P_{21}(n))P_{44}(n) + P_{41}(n)P_{24}(n) \neq 0, \quad (2.38)$$

$$(1 - P_{21}(n))(P_{43}(n) - 1) + P_{41}(n)P_{23}(n) \neq 0, \quad (2.39)$$

$$(1 - P_{43}(n))P_{24}(n) + P_{23}(n)P_{44}(n) \neq 0. \quad (2.40)$$

Then, the system (2.31) is a well-defined linear system on $C_j(n)$ which is uniquely solvable.

Proof. We first notice that system (2.31) can be rewritten as

$$\begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix} P(n)W = \begin{pmatrix} C_1(n) & 1 & C_3(n) & 0 \\ C_2(n) & 0 & C_4(n) & 1 \end{pmatrix} W = \begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix},$$

or

$$\begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix} (P(n)W - I_4) = 0. \quad (2.41)$$

Notice that

$$P(n)W - I_4 = \begin{pmatrix} P_{12}(n) - 1 & P_{11}(n) & P_{14}(n) & P_{13}(n) \\ P_{22}(n) & P_{21}(n) - 1 & P_{24}(n) & P_{23}(n) \\ P_{32}(n) & P_{31}(n) & P_{34}(n) - 1 & P_{33}(n) \\ P_{42}(n) & P_{41}(n) & P_{44}(n) & P_{43}(n) - 1 \end{pmatrix}.$$

Therefore, the conditions of this lemma mean that the second and the forth rows of $P(n)W - I_4$ are linearly independent. This, in conjunction with Remark 2.6, implies that the other two rows are linear combinations of the second and the fourth rows. In other words, there exists a *unique* collection of four numbers, $C_1(n)$, $C_2(n)$, $C_3(n)$, $C_4(n)$ such that

$$\begin{aligned} (P_{12}(n) - 1, P_{11}(n), P_{14}(n), P_{13}(n)) &= -C_1(n)(P_{22}(n), P_{21}(n) - 1, P_{24}(n), P_{23}(n)) \\ &\quad - C_3(n)(P_{42}(n), P_{41}(n), P_{44}(n), P_{43}(n) - 1), \end{aligned} \quad (2.42)$$

and

$$\begin{aligned} (P_{32}(n), P_{31}(n), P_{34}(n) - 1, P_{33}(n)) &= -C_2(n)(P_{22}(n), P_{21}(n) - 1, P_{24}(n), P_{23}(n)) \\ &\quad - C_4(n)(P_{42}(n), P_{41}(n), P_{44}(n), P_{43}(n) - 1). \end{aligned} \quad (2.43)$$

Vector equations (2.42), (2.43) are just the rows of matrix equation (2.41). ■

Investigation of the possibility of linear independence of other rows of the matrix $P(n)W - I_4$, leads us to the following observation about the uniqueness of a solution to (2.31), if one exists.

Lemma 2.8. *If the system (2.31) has a solution, it has to be unique.*

Proof. For simplicity of notation let us drop the dependence on n in $C_j(n)$ and $P(n)$, and respectively write C_j and P instead. Denote the rows of the matrix $PW - I_4$ by \mathcal{R}_j , $1 \leq j \leq 4$, and thus the system of equations (2.41) can be written as the following vector equations

$$\mathcal{R}_1 + C_1\mathcal{R}_2 + C_3\mathcal{R}_4 = 0, \quad (2.44)$$

and

$$\mathcal{R}_3 + C_2\mathcal{R}_2 + C_4\mathcal{R}_4 = 0. \quad (2.45)$$

Since $PW - I_4$ is of rank 2, necessarily, at least one of the six pairs of vectors has to be linearly independent. We consider each possibility separately.

First, let us assume that \mathcal{R}_1 and \mathcal{R}_3 are linearly independent. Note that $C_1C_4 - C_2C_3 \neq 0$, because otherwise, multiplying (2.44) by $-C_2$ and (2.45) by C_1 and adding the results yields that \mathcal{R}_1 and \mathcal{R}_3 are linearly dependent to the contrary. There exist unique constants a_1 , a_2 , b_1 , and b_2 such that

$$\mathcal{R}_2 = a_1\mathcal{R}_1 + a_2\mathcal{R}_3, \quad \text{and} \quad \mathcal{R}_4 = b_1\mathcal{R}_1 + b_2\mathcal{R}_3. \quad (2.46)$$

Also observe that, $a_1b_2 - a_2b_1 = 0$ contradicts our assumption that \mathcal{R}_1 and \mathcal{R}_3 are linearly independent. Indeed, multiplying the first member and second member of (2.46) respectively by b_2 and $-a_2$ and adding the results yields that \mathcal{R}_2 and \mathcal{R}_4 are linearly dependent, and consequently the equations (2.44) and (2.45) imply that \mathcal{R}_1 and \mathcal{R}_3 are linearly dependent, a contradiction.

Solving the system (2.44)–(2.45) for \mathcal{R}_2 and \mathcal{R}_4 and comparing with (2.46) yields

$$\begin{aligned} a_1 &= -\frac{C_4}{C_1C_4 - C_2C_3}, & a_2 &= \frac{C_3}{C_1C_4 - C_2C_3}, \\ b_1 &= \frac{C_2}{C_1C_4 - C_2C_3}, & b_2 &= -\frac{C_1}{C_1C_4 - C_2C_3}. \end{aligned}$$

Using these equalities we find $a_1b_2 - a_2b_1 = 1/(C_1C_4 - C_2C_3)$, and since $a_1b_2 - a_2b_1 \neq 0$, $C_1C_4 - C_2C_3$ is *uniquely* expressed in terms of $a_1b_2 - a_2b_1$ and hence we can write the unique solution as

$$\begin{aligned} C_1 &= -\frac{b_2}{a_1b_2 - a_2b_1}, & C_2 &= \frac{b_1}{a_1b_2 - a_2b_1}, \\ C_3 &= \frac{a_2}{a_1b_2 - a_2b_1}, & C_4 &= -\frac{a_1}{a_1b_2 - a_2b_1}. \end{aligned}$$

Next, assume that \mathcal{R}_1 and \mathcal{R}_2 are linearly independent and thus from (2.44) we necessarily have

$$C_3 \neq 0.$$

So there exist unique constants a_1 , a_2 , b_1 , and b_2 such that

$$\mathcal{R}_3 = a_1\mathcal{R}_1 + a_2\mathcal{R}_2, \quad \text{and} \quad \mathcal{R}_4 = b_1\mathcal{R}_1 + b_2\mathcal{R}_2. \quad (2.47)$$

Note that $b_1 = 0$ contradicts our assumption that \mathcal{R}_1 and \mathcal{R}_2 are linearly independent, because then $\mathcal{R}_4 = b_2\mathcal{R}_2$ and substitution into (2.44) yields $\mathcal{R}_1 = -(C_1 + b_2C_3)\mathcal{R}_2$, contrary to our assumption. Solving the system (2.44)–(2.45) for \mathcal{R}_3 and \mathcal{R}_4 and comparing with (2.47) yields the equalities:

$$\begin{aligned} a_1 &= \frac{C_4}{C_3}, & a_2 &= \frac{C_4C_1 - C_2C_3}{C_3}, \\ b_1 &= -\frac{1}{C_3}, & b_2 &= -\frac{C_1}{C_3}. \end{aligned}$$

These relationships yield the unique solution for the system (2.44)–(2.45), because C_3 is uniquely determined from b_1 (since $b_1 \neq 0$), then C_1 is uniquely determined from C_3 and b_2 , and simultaneously C_4 is uniquely determined from C_3 and a_1 , and finally C_2 is uniquely determined from C_3 , C_1 , C_4 and a_2 . Thus, we can write the unique solution as

$$\begin{aligned} C_1 &= \frac{b_2}{b_1}, & C_2 &= \frac{a_1b_2 - b_1a_2}{b_1}, \\ C_3 &= -\frac{1}{b_1}, & C_4 &= -\frac{a_1}{b_1}. \end{aligned}$$

Among the four remaining cases, the argument for each of the following three:

- \mathcal{R}_1 and \mathcal{R}_4 are linearly independent,
- \mathcal{R}_2 and \mathcal{R}_3 are linearly independent, or
- \mathcal{R}_3 and \mathcal{R}_4 are linearly independent,

is similar to the one presented above for linear independence of \mathcal{R}_1 and \mathcal{R}_2 and thus we do not provide the details here. Finally, for the case where \mathcal{R}_2 and \mathcal{R}_4 are linearly independent, we refer to Lemma 2.7 where we prove the stronger assertion that the system (2.31) is uniquely solvable. ■

Lemma 2.9. *Suppose that the solution of the \mathcal{X} -RHP exists. Then, if at least one of the conditions (2.35) through (2.40) holds, one can uniquely reconstruct the solution of the \mathcal{Y} -RHP.*

Proof. If the solution of the \mathcal{X} -RHP exists, then the expression for $P(n)$ can be found from

$$P(n) = \mathcal{X}(z; n) \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^n & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^{-n} \end{pmatrix} \Big|_{z=0}, \quad (2.48)$$

and due to our assumption, the constants $C_j(n)$ can be uniquely found according to Lemma 2.7.

Then, according to (2.29) and (2.30) we find the solution to the $\overset{\circ}{\mathcal{X}}$ -RHP as

$$\overset{\circ}{\mathcal{X}}(z; n) = \begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix} \mathcal{X}(z; n). \quad (2.49)$$

Because, if it exists, the solution of the \mathcal{X} -RHP is unique (recall Remark 2.4), we note that (2.49) is the unique solution of the $\overset{\circ}{\mathcal{X}}$ -RHP. Indeed, suppose that $\overset{\circ'}{\mathcal{X}}$ is another solution of the same $\overset{\circ}{\mathcal{X}}$ -RHP. Put

$$\mathfrak{R}'(z; n) := \overset{\circ'}{\mathcal{X}}'(z; n) \mathcal{X}^{-1}(z; n).$$

Then, as with $\mathfrak{R}(z; n)$ before, we will arrive to the conclusion that

$$\mathfrak{R}'(z; n) \equiv \begin{pmatrix} 1 & C'_1(n) & 0 & C'_3(n) \\ 0 & C'_2(n) & 1 & C'_4(n) \end{pmatrix},$$

where the constants $C'_j(n)$ satisfy the system (2.31). Since the solution of this system is unique according to Lemma 2.8, we conclude that $\mathfrak{R}'(z; n) = \mathfrak{R}(z; n)$ and, hence,

$$\overset{\circ'}{\mathcal{X}}(z; n) = \overset{\circ}{\mathcal{X}}(z; n).$$

Now, observe that the symmetry relation (2.26) implies

$$\overset{\circ}{\mathcal{X}}(z^{-1}; n) = \overset{\circ}{\mathcal{X}}(z; n) W. \quad (2.50)$$

Indeed, we have

$$\overset{\circ}{\mathcal{X}}(z^{-1}; n) = \mathfrak{R}(n) \mathcal{X}(z^{-1}; n) = \mathfrak{R}(n) P(n) W \mathcal{X}(z; n) W.$$

But, by (2.41), $\mathfrak{R}(n) P(n) W = \mathfrak{R}(n)$, therefore

$$\overset{\circ}{\mathcal{X}}(z^{-1}; n) = \mathfrak{R}(n) \mathcal{X}(z; n) W = \overset{\circ}{\mathcal{X}}(z; n) W.$$

Now, by $\overset{\circ}{\mathcal{X}}^{(j)}$, $j = 1, 2, 3, 4$ denote the columns of the 2×4 matrix $\overset{\circ}{\mathcal{X}}$. Equation (2.50) means that

$$\overset{\circ}{\mathcal{X}}^{(2)}(z; n) = \overset{\circ}{\mathcal{X}}^{(1)}(z^{-1}; n), \quad \text{and} \quad \overset{\circ}{\mathcal{X}}^{(4)}(z; n) = \overset{\circ}{\mathcal{X}}^{(3)}(z^{-1}; n).$$

In other words, the matrix valued function $\overset{\circ}{\mathcal{X}}$ can be written in the form (2.19), i.e.,

$$\overset{\circ}{\mathcal{X}}(z; n) = (\mathcal{Y}^{(1)}(z; n), \tilde{\mathcal{Y}}^{(1)}(z; n), \mathcal{Y}^{(2)}(z; n), \tilde{\mathcal{Y}}^{(2)}(z; n)),$$

with

$$\mathcal{Y}^{(1)}(z; n) := \overset{\circ}{\mathcal{X}}^{(1)}(z; n), \quad \text{and} \quad \mathcal{Y}^{(2)}(z; n) := \overset{\circ}{\mathcal{X}}^{(3)}(z; n).$$

Furthermore, the 2×2 matrix valued function

$$\mathcal{Y}(z; n) := (\mathcal{Y}^{(1)}(z; n), \mathcal{Y}^{(2)}(z; n)),$$

will be a solution of the 2×2 \mathcal{Y} -RHP. From the unique solvability of the $\overset{\circ}{\mathcal{X}}$ -RHP, it follows that this solution of the \mathcal{Y} -RHP is unique. Because otherwise, if \mathcal{Y}' is another solution of the \mathcal{Y} -RHP, via (2.19) we can construct another solution of the $\overset{\circ}{\mathcal{X}}$ -RHP on the contrary. ■

The following corollary is the direct consequence of Corollary 2.2 and Lemma 2.9.

Corollary 2.10. *Suppose that the solution of the \mathcal{X} -RHP exists for n and $n-1$, then if at least one of the conditions (2.35) through (2.40) holds also for n and $n-1$, then we have*

$$D_n \neq 0, \quad D_{n-1} \neq 0, \quad \text{and} \quad h_{n-1} \neq 0.$$

Moreover,

$$h_{n-1} = - \lim_{z \rightarrow \infty} z^{n-1} / \mathcal{Y}_{21}(z; n). \quad (2.51)$$

2.3 The primary opening of the lenses

Let us consider the contour $\Gamma := \Gamma_i \cup \mathbb{T} \cup \Gamma_o$ shown in Fig. 1. Define the function \mathcal{Z} as

$$\mathcal{Z}(z; n) := \mathcal{X}(z; n) \begin{cases} J_{\mathcal{X},i}^{-1}(z), & z \in \Omega_1, \\ J_{\mathcal{X},o}(z), & z \in \Omega_2, \\ I, & z \in \Omega_0 \cup \Omega_\infty, \end{cases}$$

where $J_{\mathcal{X},i}$ and $J_{\mathcal{X},o}$ are defined in the following factorization for the jump matrix of the \mathcal{X} -RHP, which we denote by $J_{\mathcal{X}}$:

$$\begin{aligned} J_{\mathcal{X}}(z) &:= \begin{pmatrix} 1 & 0 & \tilde{w}(z) & -\phi(z) \\ 0 & 1 & \tilde{\phi}(z) & -w(z) \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix} = \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & -w(z) \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix} \begin{pmatrix} 1 & 0 & 0 & -\phi(z) \\ 0 & 1 & \tilde{\phi}(z) & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix} \\ &\times \begin{pmatrix} 1 & 0 & \tilde{w}(z) & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix} \equiv J_{\mathcal{X},o}(z) J_{\mathcal{X},\mathbb{T}}(z) J_{\mathcal{X},i}(z). \end{aligned} \quad (2.52)$$

We remind that the symbol w , and hence \tilde{w} , are analytic in the neighborhood U_0 (cf. (1.7)) of \mathbb{T} which is supposed to include the domains Ω_1 and Ω_2 .

The function \mathcal{Z} satisfies the following Riemann–Hilbert problem:

- **RH-Z1** \mathcal{Z} is holomorphic in $\mathbb{C} \setminus (\Gamma \cup \{0\})$.

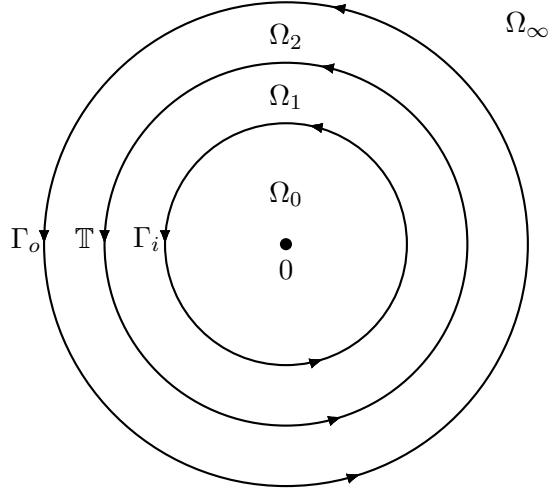


Figure 1. The jump contour Γ for the \mathcal{Z} , T and the global parametrix Riemann–Hilbert problems.

- **RH-Z2** $\mathcal{Z}_+(z; n) = \mathcal{Z}_-(z; n)J_{\mathcal{Z}}(z)$, where

$$J_{\mathcal{Z}}(z) = \begin{cases} J_{\mathcal{X}, \mathbb{T}}(z), & z \in \mathbb{T}, \\ J_{\mathcal{X}, i}(z), & z \in \Gamma_i, \\ J_{\mathcal{X}, o}(z), & z \in \Gamma_o. \end{cases}$$

- **RH-Z3** As $z \rightarrow \infty$ we have

$$\mathcal{Z}(z; n) = (I + \mathcal{O}(z^{-1})) \begin{pmatrix} z^n & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^{-n} & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}.$$

- **RH-Z4** As $z \rightarrow 0$ we have

$$\mathcal{Z}(z; n) = P(n)(I + \mathcal{O}(z)) \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^{-n} & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^n \end{pmatrix}.$$

Remark 2.11. The term “opening of the lenses” is usually used to describe situations where the jump matrix on the added contours is exponentially close to the identity matrix for large values of the parameter n . The passage $\mathbb{T} \mapsto \Gamma$, corresponding to the RH transformation $\mathcal{X} \mapsto \mathcal{Z}$, is clearly not of this type. However, our secondary opening of the lenses (the passage $\Gamma \mapsto \Gamma_S$ which corresponds to the RH transformation $T \mapsto S$) in Section 2.5 is an example of a usual opening of the lenses.

Remark 2.12. The primary opening of the lenses is essential for the progression of the RH analysis in the following sections. This is due to a technical reason that will be elaborated at the end of next section. Since the structure of jump matrices is different in Section 3, we do not have an analogous step when the Hankel symbol is supported on $[a, b]$, $0 < a < b < 1$.

2.4 Normalization of behaviours at 0 and ∞

Following the natural steps of Riemann–Hilbert analysis, we will normalize the behavior of \mathcal{Z} at 0 and ∞ ; to this end let us define

$$T(z; n) := \mathcal{Z}(z; n) \begin{cases} \begin{pmatrix} z^{-n} & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^n & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & |z| > 1, \\ \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^n & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^{-n} \end{pmatrix}, & |z| < 1. \end{cases} \quad (2.53)$$

It is very important to note that in order to have a suitable Riemann–Hilbert analysis, the normalization of behaviors at 0 and ∞ can only be carried out only after the undressing $\mathcal{X} \mapsto \mathcal{Z}$; this is due to technical reasons that will be further commented about at the end of this section. We have the following RHP for T :

- **RH-T1** T is holomorphic in $\mathbb{C} \setminus (\mathbb{T} \cup \Gamma_i \cup \Gamma_o)$.
- **RH-T2** $T_+(z; n) = T_-(z; n)J_T(z; n)$, where

$$J_T(z; n) = \begin{cases} \hat{J}(z; n), & z \in \mathbb{T}, \\ J_{\mathcal{X},i}(z), & z \in \Gamma_i, \\ J_{\mathcal{X},o}(z), & z \in \Gamma_o, \end{cases} \quad \text{where} \quad \hat{J}(z; n) = \begin{pmatrix} z^n & 0 & 0 & -\phi(z) \\ 0 & z^n & \tilde{\phi}(z) & 0 \\ 0 & 0 & z^{-n} & 0 \\ 0 & 0 & 0 & z^{-n} \end{pmatrix}, \quad (2.54)$$

and the matrices $J_{\mathcal{X},i}$ and $J_{\mathcal{X},o}$ are defined by (2.52).

- **RH-T3** As $z \rightarrow \infty$, we have $T(z; n) = (I + \mathcal{O}(z^{-1}))$.

We observe that for $z \in \mathbb{T}$, J_T can be factorized as follows

$$\begin{aligned} \hat{J}(z; n) &= \begin{pmatrix} I_2 & 0_2 \\ z^{-n}\Phi^{-1}(z) & I_2 \end{pmatrix} \begin{pmatrix} 0_2 & \Phi(z) \\ -\Phi^{-1}(z) & 0_2 \end{pmatrix} \begin{pmatrix} I_2 & 0_2 \\ z^n\Phi^{-1}(z) & I_2 \end{pmatrix} \\ &\equiv J_{T,o}(z; n) \overset{\circ}{J}(z) J_{T,i}(z; n), \end{aligned} \quad (2.55)$$

where 0_2 and I_2 are respectively 2×2 zero and identity matrices and

$$\Phi(z) = \begin{pmatrix} 0 & -\phi(z) \\ \tilde{\phi}(z) & 0 \end{pmatrix}.$$

Note that $J_{T,i}$ is exponentially close to the identity matrix for z inside of the unit circle and $J_{T,o}$ is exponentially close to the identity matrix for z outside of the unit circle.

Now we are in a position to address Remark 2.12 in the previous section. Indeed, if one normalizes the behaviors at 0 and ∞ without the undressing transformation $\mathcal{X} \mapsto \mathcal{Z}$; i.e.,

by directly defining the function \mathcal{T} as

$$\mathcal{T}(z; n) := \mathcal{X}(z; n) \begin{cases} \begin{pmatrix} z^{-n} & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^n & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & |z| > 1, \\ \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^n & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^{-n} \end{pmatrix}, & |z| < 1. \end{cases}$$

Then the jump matrix $J_{\mathcal{T}} := \mathcal{T}_-^{-1} \mathcal{T}_+$ on the unit circle would be

$$J_{\mathcal{T}}(z; n) = \begin{pmatrix} z^n & 0 & z^n \tilde{w}(z) & -\phi(z) \\ 0 & z^n & \tilde{\phi}(z) & -z^{-n} w(z) \\ 0 & 0 & z^{-n} & 0 \\ 0 & 0 & 0 & z^{-n} \end{pmatrix},$$

for which finding a factorization like (2.55) remains a challenge, mainly due to presence of the large parameter n in the 13 and 24 elements of $J_{\mathcal{T}}$. This fact justifies the necessity of the undressing step $\mathcal{X} \mapsto \mathcal{Z}$. Indeed, due to the specific matrix structure of the jump matrices $J_{\mathcal{X},i}$ and $J_{\mathcal{X},o}$ they *do not* change under the transformation (2.53).

2.5 The secondary opening of the lenses

The next Riemann–Hilbert transformation $T \mapsto S$, provides us with a problem with jump conditions on five contours where three jump matrices do not depend on n and the other two converge exponentially fast to the identity matrix as $n \rightarrow \infty$. Let us define the function S , suggested by (2.55), as

$$S(z; n) := T(z; n) \times \begin{cases} J_{T,i}^{-1}(z; n), & z \in \Omega'_1, \\ J_{T,o}(z; n), & z \in \Omega'_2, \\ I, & z \in \Omega''_1 \cup \Omega''_2 \cup \Omega_0 \cup \Omega_\infty, \end{cases}$$

where the regions Ω'_1 , Ω'_2 , Ω''_1 and Ω''_2 are shown in Fig. 2. We have the following Riemann–Hilbert problem for S

- **RH-S1** S is holomorphic in $\mathbb{C} \setminus (\mathbb{T} \cup \Gamma_i \cup \Gamma_o \cup \Gamma'_i \cup \Gamma'_o)$.
- **RH-S2** $S_+(z; n) = S_-(z; n) J_S(z; n)$, where

$$J_S(z; n) = \begin{cases} J(z), & z \in \mathbb{T}, \\ J_{T,i}(z; n), & z \in \Gamma'_i, \\ J_{T,o}(z; n), & z \in \Gamma'_o, \\ J_{\mathcal{X},i}(z), & z \in \Gamma_i, \\ J_{\mathcal{X},o}(z), & z \in \Gamma_o. \end{cases}$$

- **RH-S3** As $z \rightarrow \infty$, we have $S(z; n) = I + \mathcal{O}(z^{-1})$.

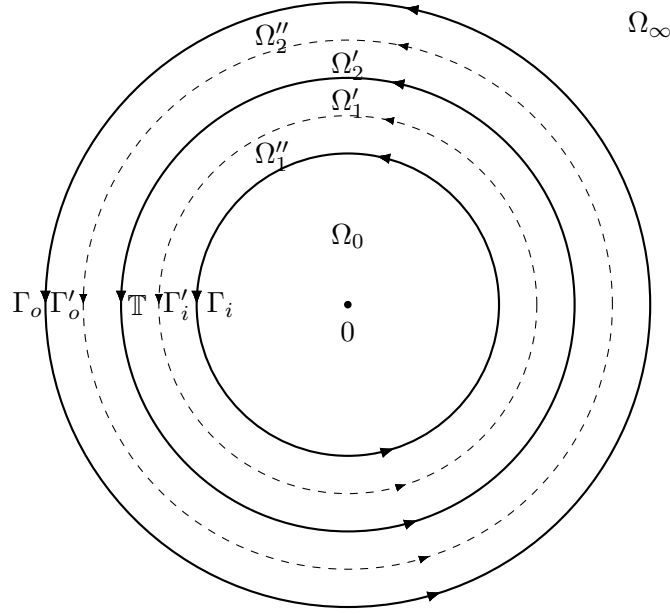


Figure 2. The jump contour Γ_S of the S -RHP.

In the usual way, we will first try to solve this Riemann–Hilbert problem by disregarding the jump matrices which depend on n , this solution is denoted by $\overset{\circ}{S}$ and will be referred to as the global parametrix. Once we construct the global parametrix, we will consider the small-norm Riemann–Hilbert problem for the ratio $R := S(\overset{\circ}{S})^{-1}$ and discuss its solvability in the forthcoming sections.

2.6 The global parametrix and the model Riemann–Hilbert problem for the pair (ϕ, w)

The S -RHP reduces to the following Riemann–Hilbert problem for the global parametrix $\overset{\circ}{S}$, when we ignore the jump matrices which are exponentially close to the identity matrix:

- **RH- $\overset{\circ}{S}1$** $\overset{\circ}{S}$ is holomorphic in $\mathbb{C} \setminus (\mathbb{T} \cup \Gamma_i \cup \Gamma_o)$.
- **RH- $\overset{\circ}{S}2$** $\overset{\circ}{S}_+(z) = \overset{\circ}{S}_-(z)J_{\overset{\circ}{S}}(z)$, where

$$J_{\overset{\circ}{S}}(z) = \begin{cases} \overset{\circ}{J}(z), & z \in \mathbb{T}, \\ J_{\mathcal{X},i}(z), & z \in \Gamma_i, \\ J_{\mathcal{X},o}(z), & z \in \Gamma_o. \end{cases}$$

- **RH- $\overset{\circ}{S}3$** As $z \rightarrow \infty$, we have $\overset{\circ}{S}(z) = I + \mathcal{O}(z^{-1})$.

And we finally *dress* the $\overset{\circ}{S}$ -RHP to obtain a model problem for the global parametrix having jumps only on the unit circle. We define the function Λ as

$$\Lambda(z) := \overset{\circ}{S}(z) \times \begin{cases} J_{\mathcal{X},i}(z), & z \in \Omega_1, \\ J_{\mathcal{X},o}^{-1}(z), & z \in \Omega_2, \\ I, & z \in \Omega_0 \cup \Omega_\infty. \end{cases} \quad (2.56)$$

Now we arrive at the following Riemann–Hilbert problem for Λ that from now on we will refer to as *the model Riemann–Hilbert problem for the pair* (ϕ, w) :

- **RH- $\Lambda 1$** Λ is holomorphic in $\mathbb{C} \setminus \mathbb{T}$.
- **RH- $\Lambda 2$** $\Lambda_+(z) = \Lambda_-(z)J_\Lambda(z)$, for $z \in \mathbb{T}$, where

$$J_\Lambda(z) = \begin{pmatrix} 0 & 0 & 0 & -\phi(z) \\ -\frac{w(z)}{\phi(z)} & 0 & \tilde{\phi}(z) - \frac{w(z)\tilde{w}(z)}{\phi(z)} & 0 \\ 0 & -\frac{1}{\tilde{\phi}(z)} & 0 & 0 \\ \frac{1}{\phi(z)} & 0 & \frac{\tilde{w}(z)}{\phi(z)} & 0 \end{pmatrix}.$$

- **RH- $\Lambda 3$** As $z \rightarrow \infty$, we have $\Lambda(z) = I + \mathcal{O}(z^{-1})$.

The conditions on w and ϕ which ensure the solvability of this model problem are not completely known and categorized at this point. We also want to stress that the appearance of the 4×4 model Λ -problem in the asymptotic analysis of the original \mathcal{X} -RHP is the crucial difference of the Toeplitz+Hankel case we consider in this work comparing to the pure Toeplitz or pure Hankel or Toeplitz+Hankel with the same symbols cases. Indeed, even if the pair (ϕ, w) is such that the Λ -RHP is solvable it does not mean that it is *explicitly* solvable. Hence, one should not expect the closed form of the asymptotic answer in the case of the generic pair (ϕ, w) .⁴ However, in Section 4 we will present a detailed analysis of this model problem for a specific family of pairs (ϕ, w) within the broader class of Toeplitz and Hankel weights considered by E. Basor and T. Ehrhardt in [2] for which the model Λ -problem is explicitly solvable.

Remarkably, we arrive at the *same* model Λ -Riemann–Hilbert problem, if we start with a Hankel weight supported on the interval $[a, b]$, $0 < a < b < 1$. This will be shown in the next section.

3 Toeplitz+Hankel determinants: Hankel weight supported on the interval $[a, b]$, $0 < a < b < 1$

In this section we consider the determinant (1.3) where w_k exist and are given by (1.2), with $I = [a, b]$, $0 < a < b < 1$. We further assume that w does not have Fisher–Hartwig singularities (see [10, 11, 19], or [20] for instances of Fisher–Hartwig singularities on the real line). Let us again assume that the symbol ϕ is of Szegő-type. The Riemann–Hilbert approach outlined in this section can be naturally extended to the three other cases: *i*) $-1 < a < b < 0$, *ii*) $-\infty < a < b < -1$, and *iii*) $1 < a < b < \infty$. We consider the system of orthogonal polynomials $\{P_n(z)\}$, $\deg P_n(z) = n$, satisfying the following orthogonality conditions

$$\int_a^b P_n(x)x^{k+s}w(x)dx + \int_{\mathbb{T}} P_n(z)z^{-k-r}\phi(z)\frac{dz}{2\pi iz} = h_n\delta_{n,k}, \quad k = 0, 1, \dots, n. \quad (3.1)$$

One can write a determinantal formula for P_n like in (2.2), which yields

$$h_n = \frac{D_{n+1}(\phi, w; r, s)}{D_n(\phi, w; r, s)}.$$

⁴In this respect the Toeplitz+Hankel determinants with generic symbols are similar to the block Toeplitz determinants, where the explicit answers can be obtained only in two cases: (a) the Fourier expansion of the corresponding matrix symbol is one side truncated or (b) one can produce an explicit Wiener–Hopf factorization of the symbol.

By similar considerations as those mentioned in Section 2, the orthogonal polynomials P_n exist and are unique if $D_n \neq 0$. Now, assuming that $D_n, D_{n-1} \neq 0$, we consider the function Y defined as

$$Y(z; n) = \begin{pmatrix} P_n(z) & \int_a^b \frac{P_n(x)x^s w(x)}{x-z} dx + \int_{\mathbb{T}} \frac{\tilde{\phi}(\xi)\xi^r \tilde{P}_n(\xi)}{\xi-z} \frac{d\xi}{2\pi i \xi} \\ -\frac{P_{n-1}(z)}{h_{n-1}} & -\frac{1}{h_{n-1}} \left\{ \int_a^b \frac{P_{n-1}(x)x^s w(x)}{x-z} dx + \int_{\mathbb{T}} \frac{\tilde{\phi}(\xi)\xi^r \tilde{P}_{n-1}(\xi)}{\xi-z} \frac{d\xi}{2\pi i \xi} \right\} \end{pmatrix}, \quad (3.2)$$

built from the orthogonal polynomials P_n satisfying (3.1). Consider the following Riemann–Hilbert problem for finding the 2×2 matrix Y satisfying

- **RH-Y1** Y is holomorphic in $\mathbb{C} \setminus (\mathbb{T} \cup [a, b])$.
- **RH-Y2** For $z \in \mathbb{T}$ we have

$$Y_+^{(1)}(z; n) = Y_-^{(1)}(z; n),$$

and

$$Y_+^{(2)}(z; n) = Y_-^{(2)}(z; n) + z^{r-1} \tilde{\phi}(z) Y_-^{(1)}(z^{-1}; n).$$

- **RH-Y3** For $x \in (a, b)$ we have

$$Y_+^{(1)}(x; n) = Y_-^{(1)}(x; n),$$

and

$$Y_+^{(2)}(x; n) = Y_-^{(2)}(x; n) + 2\pi i x^s w(x) Y_-^{(1)}(x; n).$$

- **RH-Y4** As $z \rightarrow \infty$

$$Y(z; n) = \left(I + O\left(\frac{1}{z}\right) \right) z^{n\sigma_3} = \begin{pmatrix} z^n + \mathcal{O}(z^{n-1}) & \mathcal{O}(z^{-n-1}) \\ \mathcal{O}(z^{n-1}) & z^{-n} + \mathcal{O}(z^{-n-1}) \end{pmatrix},$$

where $Y^{(1)}$ and $Y^{(2)}$ are the first and second columns of Y , respectively. We have the analogue of Theorem 2.1 here as well.

Theorem 3.1. *The following statements are true.*

1. Suppose that $D_n, D_{n-1} \neq 0$. Then, the Riemann–Hilbert problem **RH-Y1** through **RH-Y4** is uniquely solvable and its solution \mathcal{Y} is defined by (3.2). Moreover,

$$h_{n-1} = -\lim_{z \rightarrow \infty} z^{n-1} / Y_{21}(z; n).$$

2. Suppose that the Riemann–Hilbert problem **RH-Y1** through **RH-Y4** has a unique solution. Then $D_n \neq 0$, $\text{rank}(T_{n-1}[\phi; r] + H_{n-1}[w; s]) \geq n-2$, and $P_n(z) = Y_{11}(z; n)$.
3. Suppose that the Riemann–Hilbert problem **RH-Y1** through **RH-Y4** has a unique solution. Suppose also that

$$\lim_{z \rightarrow \infty} Y_{21}(z; n) z^{-n+1} \neq 0.$$

Then, as before, $D_n \neq 0$, $P_n(z) = Y_{11}(z; n)$, and, in addition,

$$D_{n-1} \neq 0, \quad h_{n-1} = -\lim_{z \rightarrow \infty} \mathcal{Y}_{21}^{-1}(z; n) z^{n-1}, \quad P_{n-1}(z) = -h_{n-1} Y_{21}(z; n).$$

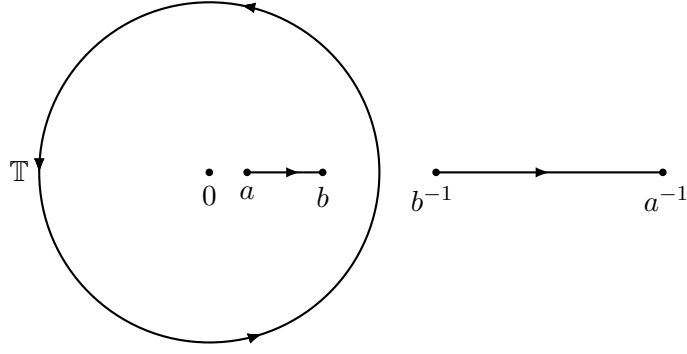


Figure 3. The jump contour Σ .

We omit the proof here as it is similar to the proof of Theorem 2.1.

Corollary 3.2. *Suppose that the Y -RH problem has a unique solution for n and $n - 1$. Then*

$$D_n \neq 0, \quad D_{n-1} \neq 0, \quad \text{and} \quad h_{n-1} \neq 0.$$

Moreover,

$$h_{n-1} = - \lim_{z \rightarrow \infty} z^{n-1} / Y_{21}(z; n).$$

3.1 The associated 2×4 and 4×4 Riemann–Hilbert problems

The formulation of the 2×4 and 4×4 Riemann–Hilbert problems are very similar to those of Section 2.1, however there are minor differences that convinces us to practice clarity in our exposition. Let us consider the following 2×4 matrix function, constructed from the columns of Y given by (3.2):

$$\mathring{X}(z; n) := (Y^{(1)}(z; n), \tilde{Y}^{(1)}(z; n), Y^{(2)}(z; n), \tilde{Y}^{(2)}(z; n)). \quad (3.3)$$

Let us define $\Sigma := \mathbb{T} \cup [a, b] \cup [b^{-1}, a^{-1}]$, and $\Sigma' := \Sigma \setminus \{a, b, b^{-1}, a^{-1}\}$. $\mathring{X}(z; n)$ satisfies the following Riemann–Hilbert problem

- **RH- $\mathring{X}1$** \mathring{X} is analytic in $\mathbb{C} \setminus (\Sigma \cup \{0\})$,
- **RH- $\mathring{X}2$** For $z \in \Sigma'$, we have $\mathring{X}_+(z; n) = \mathring{X}_-(z; n) J_{\mathring{X}}(z)$, where

$$J_{\mathring{X}}(z) = \begin{cases} \begin{pmatrix} 1 & 0 & 0 & -z^{-r+1}\phi(z) \\ 0 & 1 & z^{r-1}\tilde{\phi}(z) & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & z \in \mathbb{T}, \\ \begin{pmatrix} 1 & 0 & 2\pi i x^s w(x) & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & z \equiv x \in (a, b), \\ \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & -2\pi i x^{-s} \tilde{w}(x) \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & z \equiv x \in (b^{-1}, a^{-1}). \end{cases}$$

- **RH- $\overset{\circ}{X}3$** As $z \rightarrow \infty$

$$\overset{\circ}{X}(z; n) = \begin{pmatrix} 1 + O(z^{-1}) & E_1(n) + O(z^{-1}) & O(z^{-1}) & E_3(n) + O(z^{-1}) \\ O(z^{-1}) & E_2(n) + O(z^{-1}) & 1 + O(z^{-1}) & E_4(n) + O(z^{-1}) \end{pmatrix} \\ \times \begin{pmatrix} z^n & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^{-n} & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}.$$

- **RH- $\overset{\circ}{X}4$** As $z \rightarrow 0$

$$\overset{\circ}{X}(z; n) = \begin{pmatrix} E_1(n) + O(z) & 1 + O(z) & E_3(n) + O(z) & O(z) \\ E_2(n) + O(z) & O(z) & E_4(n) + O(z) & 1 + O(z) \end{pmatrix} \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^{-n} & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^n \end{pmatrix},$$

where

$$\begin{aligned} E_1(n) &= Y_{11}(0; n), & E_3(n) &= Y_{12}(0; n), \\ E_2(n) &= Y_{21}(0; n), & E_4(n) &= Y_{22}(0; n). \end{aligned}$$

It is straightforward to check that $\overset{\circ}{X}$ given by (3.3) and (3.2) satisfies the Riemann–Hilbert problem **RH- $\overset{\circ}{X}1$** through **RH- $\overset{\circ}{X}4$** .

One of the differences between the case when the Hankel symbol is supported on the unit circle versus the case when it is supported on the interval $[a, b]$, is discussed in the following remark about the values of offsets that can be handled without much difficulty in each case.

Remark 3.3. Let u be defined by

$$u(z) := z \int_a^b \frac{t^{s-1} w(t)}{t - z} dt. \quad (3.4)$$

When the Hankel symbol is supported on the interval, the natural progression of the Riemann–Hilbert analysis with general offset values $r, s \in \mathbb{Z}$, finally requires us to construct Szegő functions for the functions $f_1(z) = z^{r-1} \phi(z)$ and $f_2(z) = \tilde{u}(z)$ (Compare with Remark 2.3). Note that the function u given by (3.4) (and \tilde{u}) has no winding number for all $s \in \mathbb{Z}$. Therefore, in this work it seems natural for us to focus on the determinants of the type $D_n(\phi, w; 1, s)$, in view of the points mentioned above and in Remark 2.3.

Similar to our approach in Section 2.1, we introduce the following Riemann–Hilbert problem of finding the 4×4 matrix function X satisfying:

- **RH-X1** X is analytic in $\mathbb{C} \setminus (\Sigma \cup \{0\})$.

- **RH-X2** For $z \in \Sigma'$, we have $X_+(z; n) = X_-(z; n)J_X(z)$, where

$$J_X(z) = \begin{cases} \begin{pmatrix} 1 & 0 & 0 & -\phi(z) \\ 0 & 1 & \tilde{\phi}(z) & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & z \in \mathbb{T}, \\ \begin{pmatrix} 1 & 0 & 2\pi i x^s w(x) & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & z \equiv x \in (a, b), \\ \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & -2\pi i x^{-s} \tilde{w}(x) \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & z \equiv x \in (b^{-1}, a^{-1}), \end{cases} \quad (3.5)$$

that is, J_X is exactly equal to $J_{\overset{\circ}{X}}$ when $r = 1$.

- **RH-X3** As $z \rightarrow \infty$

$$X(z; n) = (I + \mathcal{O}(z^{-1})) \begin{pmatrix} z^n & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^{-n} & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}.$$

- **RH-X4** As $z \rightarrow 0$

$$X(z; n) = Q(n)(I + \mathcal{O}(z)) \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^{-n} & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^n \end{pmatrix}, \quad (3.6)$$

where we emphasize that the matrix factor $Q(n)$ in (3.6) is not a priori prescribed. Now we are going to briefly mention some facts about this Riemann–Hilbert problem which are similar to those of the \mathcal{X} -RHP. Using the usual Liouville theorem-based arguments one can easily show that the solution of X -RHP is unique, if it exists. Also, without much difficulty one can show that the function $WP^{-1}(n)X(z^{-1}; n)W$ is also a solution of the X -RHP, and thus due to the uniqueness of the solution, we get the symmetry relation

$$WP^{-1}(n)\mathcal{X}(z^{-1}; n)W = \mathcal{X}(z; n), \quad (3.7)$$

where W is given by (2.27). Equation (3.7) yields the following symmetry equation for $Q(n)$,

$$Q(n) = WQ^{-1}(n)W,$$

or,

$$(WQ(n))^2 = (Q(n)W)^2 = I_4.$$

Exact similar argument used in Remark 2.6 proves that the matrix $Q(n)W - I_4$ has rank 2. Here we also have the relationship between $\overset{\circ}{X}$ and X given by

$$\overset{\circ}{X}(z; n) = \begin{pmatrix} 1 & E_1(n) & 0 & E_3(n) \\ 0 & E_2(n) & 1 & E_4(n) \end{pmatrix} X(z; n),$$

and moreover,

$$\begin{pmatrix} 1 & E_1(n) & 0 & E_3(n) \\ 0 & E_2(n) & 1 & E_4(n) \end{pmatrix} = \begin{pmatrix} E_1(n) & 1 & E_3(n) & 0 \\ E_2(n) & 0 & E_4(n) & 1 \end{pmatrix} Q^{-1}(n). \quad (3.8)$$

Since this is exactly the system (2.31), where $C_j(n)$ and $P(n)$ are respectively replaced by $E_j(n)$ and $Q(n)$, $1 \leq j \leq 4$, and because $P(n)$ and $Q(n)$ enjoy the same symmetry and rank properties, we readily have the following statements, whose counterparts are already proven in Section 2.

Lemma 3.4. *Let $Q_{jk}(n)$, $1 \leq j, k \leq 4$, be the entries of the matrix $Q(n)$. Assume that at least one of the following six inequalities is true,*

$$Q_{22}(n)Q_{44}(n) - Q_{42}(n)Q_{24}(n) \neq 0, \quad (3.9)$$

$$(1 - Q_{21}(n))Q_{42}(n) + Q_{22}(n)Q_{41}(n) \neq 0, \quad (3.10)$$

$$(1 - Q_{43}(n))Q_{22}(n) + Q_{23}(n)Q_{42}(n) \neq 0, \quad (3.11)$$

$$(1 - Q_{21}(n))Q_{44}(n) + Q_{41}(n)Q_{24}(n) \neq 0, \quad (3.12)$$

$$(1 - Q_{21}(n))(Q_{43}(n) - 1) + Q_{41}(n)Q_{23}(n) \neq 0, \quad (3.13)$$

$$(1 - Q_{43}(n))Q_{24}(n) + Q_{23}(n)Q_{44}(n) \neq 0. \quad (3.14)$$

Then, the system (3.8) is a well-defined linear system on $E_j(n)$ which is uniquely solvable.

Lemma 3.5. *If the system (3.8) has a solution, it has to be unique.*

Lemma 3.6. *Suppose that the solution of the X-RHP exists. Then, if at least one of the conditions (3.9) through (3.14) holds, one can uniquely reconstruct the solution of the Y-RHP.*

Corollary 3.7. *Suppose that the solution of the X-RHP exists for n and $n - 1$, then if at least one of the conditions (3.9) through (3.14) also holds for n and $n - 1$, then we have*

$$D_n \neq 0, \quad D_{n-1} \neq 0, \quad \text{and} \quad h_{n-1} \neq 0,$$

where h_{n-1} can be reconstructed from the RHP data as

$$h_{n-1} = - \lim_{z \rightarrow \infty} z^{n-1} / Y_{21}(z; n).$$

3.2 Normalization of behaviors at 0 and ∞

Unlike the situation in Section 2 where we had to make the transformation $\mathcal{X} \mapsto \mathcal{Z}$ before normalization of behaviors at zero and infinity, when the Hankel symbol is supported on the interval $[a, b]$ we can immediately normalize the asymptotic behaviors at 0 and infinity, due to the desired structure of jump matrices. Indeed, it is natural to define

$$T(z; n) := X(z; n) \begin{cases} \begin{pmatrix} z^{-n} & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & z^n & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & |z| > 1, \\ \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^n & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^{-n} \end{pmatrix}, & |z| < 1. \end{cases} \quad (3.15)$$

The function T satisfies the following Riemann–Hilbert problem:

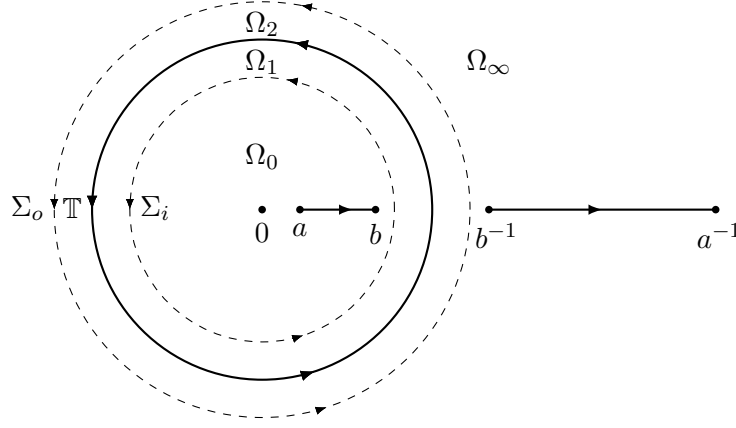


Figure 4. The jump contour $\Sigma_S \equiv \Sigma \cup \Sigma_o \cup \Sigma_i$ of the S -RHP.

- **RH-T1** T is holomorphic in $\mathbb{C} \setminus \Sigma$.
- **RH-T2** For $z \in \Sigma'$, we have $T_+(z; n) = T_-(z; n)J_T(z; n)$, where

$$J_T(z; n) = \begin{cases} \hat{J}(z; n), & z \in \mathbb{T}, \\ J_X(z), & z \in (a, b) \cup (b^{-1}, a^{-1}). \end{cases}$$

We recall that \hat{J} is given by (2.54) and the matrices J_X for $z \in (a, b)$ and $z \in (b^{-1}, a^{-1})$ are given by (3.5).

- **RH-T3** As $z \rightarrow \infty$, we have $T(z; n) = (I + \mathcal{O}(z^{-1}))$.

We bring the reader's attention to the fact that the transformation (3.15) does not change the jump matrices J_X .

3.3 Opening of the lenses

Using (2.55), we open the lenses off the unit circle as shown in the Fig. 4 and we define

$$S(z; n) := T(z; n) \times \begin{cases} J_{T,i}^{-1}(z; n), & z \in \Omega_1, \\ J_{T,o}(z; n), & z \in \Omega_2, \\ I, & z \in \mathbb{C} \setminus (\overline{\Omega_1} \cup \overline{\Omega_2} \cup [a, b] \cup [b^{-1}, a^{-1}]), \end{cases}$$

where $J_{T,i}$ and $J_{T,o}$ are defined in (2.55). Let $\Sigma_S \equiv \Sigma \cup \Sigma_o \cup \Sigma_i$ and $\Sigma'_S \equiv \Sigma' \cup \Sigma_o \cup \Sigma_i$ (see Fig. 4). It is straightforward to check that S satisfies the following Riemann–Hilbert problem

- **RH-S1** S is holomorphic in $\mathbb{C} \setminus \Sigma_S$,
- **RH-S2** For $z \in \Sigma'_S$ we have $S_+(z; n) = S_-(z; n)J_S(z; n)$, where

$$J_S(z; n) = \begin{cases} \overset{\circ}{J}(z), & z \in \mathbb{T}, \\ J_{T,i}(z; n), & z \in \Sigma_i, \\ J_{T,o}(z; n), & z \in \Sigma_o, \\ J_X(z), & z \in (a, b) \cup (b^{-1}, a^{-1}), \end{cases}$$

where these matrices are defined in (2.55) and (3.5).

- **RH-S3** As $z \rightarrow \infty$, we have $S(z; n) = I + \mathcal{O}(z^{-1})$.

3.4 The global parametrix and a model Riemann–Hilbert problem

Let us consider the following Riemann–Hilbert for \mathring{S} , or the global parametrix, which is expected to be a good approximation to S for large parameter n . This RHP is simply obtained from the S -RHP by ignoring the jumps on Σ_i and Σ_o :

- **RH- $\mathring{S}1$** \mathring{S} is holomorphic in $\mathbb{C} \setminus \Sigma$ (see Fig. 3).
- **RH- $\mathring{S}2$** For $z \in \Sigma'$, we have $\mathring{S}_+(z) = \mathring{S}_-(z)J_{\mathring{S}}(z)$, where

$$J_{\mathring{S}}(z) = \begin{cases} \mathring{J}(z), & z \in \mathbb{T}, \\ J_X(z), & z \in (a, b) \cup (b^{-1}, a^{-1}). \end{cases}$$

- **RH- $\mathring{S}3$** As $z \rightarrow \infty$, we have $\mathring{S}(z) = I + \mathcal{O}(z^{-1})$.

Let us recall the function u defined in Remark 3.3:

$$u(z) = z \int_a^b \frac{t^{s-1}w(t)}{t-z} dt.$$

The Plemelj–Sokhotskii formula implies that

$$\begin{aligned} u_+(x) - u_-(x) &= 2\pi i x^s w(x), & x \in (a, b), \\ \tilde{u}_+(x) - \tilde{u}_-(x) &= -2\pi i x^{-s} \tilde{w}(x), & x \in (b^{-1}, a^{-1}). \end{aligned}$$

Put

$$\Theta(z) := \mathring{S}(z) \begin{cases} \begin{pmatrix} 1 & 0 & -u(z) & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & |z| < 1, \\ \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & -\tilde{u}(z) \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}, & |z| > 1. \end{cases}$$

It can be checked that Θ does not have jumps on the intervals (a, b) and (b^{-1}, a^{-1}) . We have arrived at the following *model Riemann–Hilbert problem* on the unit circle:

- **RH- $\Theta1$** Θ is holomorphic in $\mathbb{C} \setminus \mathbb{T}$.
- **RH- $\Theta2$** $\Theta_+(z) = \Theta_-(z)J_{\Theta}(z)$, for $z \in \mathbb{T}$, where

$$J_{\Theta}(z) = \begin{pmatrix} 0 & 0 & 0 & -\phi(z) \\ \frac{\tilde{u}(z)}{\phi(z)} & 0 & \tilde{\phi}(z) - \frac{u(z)\tilde{u}(z)}{\phi(z)} & 0 \\ 0 & -\frac{1}{\tilde{\phi}(z)} & 0 & 0 \\ \frac{1}{\phi(z)} & 0 & -\frac{u(z)}{\phi(z)} & 0 \end{pmatrix}.$$

- **RH-Θ3** As $z \rightarrow \infty$, we have $\Theta(z) = I + \mathcal{O}(z^{-1})$,

where, in **RH-Θ3** we have used the fact that $\tilde{u}(\infty) = u(0) = 0$.

Remark 3.8. Recalling Section 2.6, we note that this is exactly the model Riemann–Hilbert problem for the pair $(\phi, -\tilde{u})$. Hence, it can be concluded that the study of the Toeplitz+Hankel determinants both when the Hankel symbol is supported on the unit circle and also when it is supported on the interval $[a, b]$, reduces to the study of the model Riemann–Hilbert problem **RH-Λ1** through **RH-Λ3**.

For a specific class of symbols ϕ and w , we will present the solution to the model Riemann–Hilbert problem **RH-Λ1** through **RH-Λ3** for the pair (ϕ, w) in the next section.

4 Analysis of the model problem and a solvable pair

As mentioned before, it is an ambitious task to classify all the pairs (ϕ, w) for which the Λ -model Riemann–Hilbert problem is solvable. However, it is reasonable to start our analysis with the class of symbols (1.4) considered in [2]. Since in our work the symbols are not assumed to be of the Fisher–Hartwig type (which needs a more delicate treatment, see Section 5.4), we should still expect that the model Riemann–Hilbert problem be solvable for the class of symbols (1.4) when there is no Fisher–Hartwig singularity ($a_0(z) = b_0(z) \equiv 1$). Indeed this is the case as will be elaborated in this section. As commented in the beginning of Section 2.1, asymptotics of $D_n(\phi, d\phi; r, s)$, for general r and s requires a more delicate approach (see Section 5.2) and we do not discuss the details here. So let us consider $D_n(\phi, d\phi; 1, 1)$, where d is of Szegő-type and further satisfies the condition $d(z)\tilde{d}(z) = 1$ on the unit circle. For instance, a class of functions satisfying these conditions is given by

$$d(z) = \prod_{j=1}^m d_j(z), \quad d_j(z) = \pm \left(\frac{z - b_j}{z - a_j} \right)^{\alpha_j} \left(\frac{a_j z - 1}{b_j z - 1} \right)^{\alpha_j}, \quad (4.1)$$

where $\alpha_j \in \mathbb{C}$, all factors are defined by their principal branch, and

$$0 < a_1 < b_1 < a_2 < b_2 < \cdots < a_m < b_m < 1.$$

Note that a similar construction can be found for $-1 < b_m < a_m < \cdots < b_1 < a_1 < 0$, and thus a larger class of functions can be found from multiplying functions of the first class with those of the second class. Although we have a class of functions satisfying the required properties, a complete categorization of functions satisfying the four required properties for d is yet to be found. We emphasize that the conditions $d(\pm 1) = 1$ required in [2] do not play a role in the Riemann–Hilbert analysis. Indeed, for d as defined in (4.1) one can check that $d(\pm 1) = (-1)^{\epsilon_0}$, where ϵ_0 is the number of the d_j -factors in whose definition the sign “−” is taken. So in this sense we are considering functions d which are slightly more general than those considered in [2]. At the same time, we have our technical assumption of analyticity of the symbols in a neighborhood of the unit circle which is not needed in the analysis of [2].

Note that the condition $d\tilde{d} = 1$ on the unit circle renders the 23-element of the jump matrix J_Λ zero; indeed

$$J_{\Lambda, 23}(z) = \tilde{\phi}(z) - \frac{w(z)\tilde{w}(z)}{\phi(z)} = \tilde{\phi}(z)(1 - d(z)\tilde{d}(z)) = 0.$$

Hence, for the particular choices made above, the jump matrix G_Λ reduces to

$$J_\Lambda(z) = \begin{pmatrix} 0 & 0 & 0 & -\phi(z) \\ -d(z) & 0 & 0 & 0 \\ 0 & -\frac{1}{\tilde{\phi}(z)} & 0 & 0 \\ \frac{1}{\phi(z)} & 0 & \frac{\tilde{w}(z)}{\phi(z)} & 0 \end{pmatrix}. \quad (4.2)$$

In order to factorize J_Λ , let us first consider the following Szegő functions

$$\alpha(z) = \exp \left[\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(\phi(\tau))}{\tau - z} d\tau \right], \quad \beta(z) = \exp \left[\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(d(\tau))}{\tau - z} d\tau \right]. \quad (4.3)$$

By Plemelj–Sokhotskii formula α , β , $\tilde{\alpha}$ and $\tilde{\beta}$ satisfy the following jump conditions on the unit circle:

$$\begin{aligned} \alpha_+(z) &= \alpha_-(z)\phi(z), & \beta_+(z) &= \beta_-(z)d(z), \\ \tilde{\alpha}_-(z) &= \tilde{\alpha}_+(z)\tilde{\phi}(z), & \tilde{\beta}_-(z) &= \tilde{\beta}_+(z)\tilde{d}(z). \end{aligned} \quad (4.4)$$

It turns out that knowing the value of $\beta(0)$ is crucial for finding an asymptotic expression for h_n (see Section 4.2) and the condition $d\tilde{d} = 1$ on the unit circle allows us to evaluate $\beta(0)$ easily. Indeed

$$\int_{\mathbb{T}} \ln(d(\tau)) \frac{d\tau}{\tau} = \int_{\mathbb{T}} \ln(\tilde{d}(\tau)) \frac{d\tau}{\tau} = \int_{\mathbb{T}} \ln(d^{-1}(\tau)) \frac{d\tau}{\tau} = - \int_{\mathbb{T}} \ln(d(\tau)) \frac{d\tau}{\tau}.$$

Thus

$$\int_{\mathbb{T}} \ln(d(\tau)) \frac{d\tau}{\tau} = 0, \quad \text{and therefore,} \quad \beta(0) = 1. \quad (4.5)$$

Next, we show that $\beta = \tilde{\beta}$. Note that

$$\begin{aligned} \tilde{\beta}(z) &= \exp \left[\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(d(\tau))}{\tau - z^{-1}} d\tau \right] = \exp \left[-\frac{z}{2\pi i} \int_{\mathbb{T}} \frac{\ln(d(\tau))}{\tau^{-1} - z} \frac{d\tau}{\tau} \right] \\ &= \exp \left[-\frac{z}{2\pi i} \int_{\mathbb{T}} \frac{\ln(\tilde{d}(\tau))}{\tau - z} \frac{d\tau}{\tau} \right] = \exp \left[\frac{z}{2\pi i} \int_{\mathbb{T}} \frac{\ln(d(\tau))}{\tau - z} \frac{d\tau}{\tau} \right], \end{aligned}$$

where we have again used the fact that $d\tilde{d} \equiv 1$ on the unit circle. Using

$$\frac{1}{(\tau - z)\tau} = \frac{z^{-1}}{\tau - z} - \frac{z^{-1}}{\tau},$$

we can write the last expression for $\tilde{\beta}$ as

$$\tilde{\beta}(z) = \exp \left[\frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(d(\tau))}{\tau - z} d\tau - \frac{1}{2\pi i} \int_{\mathbb{T}} \frac{\ln(d(\tau))}{\tau} d\tau \right] = \frac{\beta(z)}{\beta(0)} = \beta(z),$$

by (4.5). To show that $\beta = \tilde{\beta}$ one could also argue that they both solve *the same* scalar RHP which has a unique solution. We also note that $\alpha(z), \beta(z) = 1 + \mathcal{O}(z^{-1})$, and $\tilde{\alpha}(z) =$

$\alpha(0)(1 + \mathcal{O}(z^{-1}))$ as $z \rightarrow \infty$. Now we can write the solution of the Λ -RHP (in the case $d\tilde{d} \equiv 1$ on \mathbb{T}) as

$$\Lambda(z) = \Lambda_\infty^{-1} \begin{pmatrix} 1 & 0 & 0 & 0 \\ \mathcal{C}_\rho(z) & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix} \times \begin{cases} \begin{pmatrix} -\beta(z) & 0 & 0 & 0 \\ 0 & 0 & \frac{1}{\tilde{\alpha}(z)\beta(z)\alpha(z)} & 0 \\ 0 & -\tilde{\alpha}(z) & 0 & 0 \\ 0 & 0 & 0 & -\alpha(z) \end{pmatrix}, & |z| < 1, \\ \begin{pmatrix} 0 & \beta(z) & 0 & 0 \\ 0 & 0 & 0 & \frac{1}{\beta(z)\tilde{\alpha}(z)\alpha(z)} \\ 0 & 0 & \tilde{\alpha}(z) & 0 \\ \alpha(z) & 0 & 0 & 0 \end{pmatrix}, & |z| > 1, \end{cases} \quad (4.6)$$

where $\mathcal{C}_f(z)$ is the Cauchy-transform of $f(z)$:

$$\mathcal{C}_f(z) = \frac{1}{2\pi i} \int_{\mathbb{T}} \frac{f(\tau)}{\tau - z} d\tau,$$

and

$$\Lambda_\infty^{-1} = \begin{pmatrix} 0 & 0 & 0 & 1 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & \frac{1}{\alpha(0)} & 0 \\ 0 & \alpha(0) & 0 & 0 \end{pmatrix}, \quad \rho(z) = -\frac{1}{\beta_-(z)\beta_+(z)\tilde{\alpha}_-(z)\alpha_+(z)}. \quad (4.7)$$

Using (4.4), the Plemelj–Sokhotskii formula and general properties of the Cauchy integral, it can be checked that Λ given by (4.6) satisfies the Λ -RHP.

4.1 The small-norm Riemann–Hilbert problem associated to $D_n(\phi, d\phi, 1, 1)$

Let us consider

$$R(z; n) := S(z; n) \overset{\circ}{S}(z)^{-1}.$$

This function clearly has no jumps on Γ_i , Γ_o and \mathbb{T} , since S and $\overset{\circ}{S}$ have the same jumps on these contours. Thus, R satisfies the following small-norm Riemann–Hilbert problem

- **RH-R1** R is holomorphic in $\mathbb{C} \setminus \Gamma_R$.
- **RH-R2** $R_+(z; n) = R_-(z; n)J_R(z; n)$, for $z \in \Gamma_R$.
- **RH-R3** As $z \rightarrow \infty$, $R(z; n) = I + \mathcal{O}(z^{-1})$,

where $\Gamma_R := \Gamma'_i \cup \Gamma'_o$, and J_R is given by

$$J_R(z; n) = \overset{\circ}{S}(z)J_S(z; n)\overset{\circ}{S}(z)^{-1} = \begin{cases} \overset{\circ}{S}(z)J_{T,i}(z; n)\overset{\circ}{S}(z)^{-1}, & z \in \Gamma'_i, \\ \overset{\circ}{S}(z)J_{T,o}(z; n)\overset{\circ}{S}(z)^{-1}, & z \in \Gamma'_o. \end{cases}$$

Using (4.6), (4.7), (2.56), (2.55) and (2.52) we find

$$J_R(z; n) - I = \begin{cases} z^n \begin{pmatrix} 0 & g_{12}(z) & 0 & g_{14}(z) \\ 0 & 0 & g_{23}(z) & 0 \\ 0 & 0 & 0 & 0 \\ 0 & 0 & g_{43}(z) & 0 \end{pmatrix}, & z \in \Gamma'_i, \\ z^{-n} \begin{pmatrix} 0 & 0 & 0 & 0 \\ g_{21}(z) & 0 & 0 & 0 \\ 0 & g_{32}(z) & 0 & g_{34}(z) \\ g_{41}(z) & 0 & 0 & 0 \end{pmatrix}, & z \in \Gamma'_o, \end{cases} \quad (4.8)$$

where

$$\begin{aligned} g_{12}(z) &= -\frac{\alpha(z)}{\phi(z)\beta(z)} - \frac{\tilde{w}(z)C_\rho(z)}{\phi(z)\beta(z)\tilde{\alpha}(z)}, & g_{14}(z) &= \frac{\tilde{w}(z)}{\phi(z)\beta(z)\tilde{\alpha}(z)\alpha(0)}, \\ g_{23}(z) &= -\frac{\alpha(0)\tilde{w}(z)\beta(z)}{\tilde{\phi}(z)\tilde{\alpha}(z)}, & g_{43}(z) &= -\alpha^2(0)\left(\frac{\alpha(z)\beta(z)}{\tilde{\phi}(z)} + \frac{\beta(z)\tilde{w}(z)C_\rho(z)}{\tilde{\alpha}(z)\tilde{\phi}(z)}\right), \\ g_{21}(z) &= \frac{w(z)\beta(z)}{\phi(z)\alpha(z)}, & g_{32}(z) &= -\frac{1}{\alpha(0)\tilde{\phi}(z)}\left(\frac{\tilde{\alpha}(z)}{\beta(z)} - w(z)\tilde{\alpha}^2(z)\beta(z)\alpha(z)C_\rho(z)\right), \\ g_{34}(z) &= \frac{w(z)\tilde{\alpha}^2(z)\beta(z)\alpha(z)}{\tilde{\phi}(z)\alpha^2(0)}, & g_{41}(z) &= -\frac{\alpha(0)}{\phi(z)}\left(\frac{1}{\tilde{\alpha}(z)\beta(z)\alpha^2(z)} - \frac{w(z)\beta(z)C_\rho(z)}{\alpha(z)}\right). \end{aligned}$$

From (4.8) it follows that the jump matrix J_R satisfies on Γ_R the small-norm estimate,

$$\|J_R - I\|_{L_2 \cap L_\infty} \leq Ce^{-cn}, \quad (4.9)$$

for some positive C and $c = -\log r_1$, where r_1 is any number satisfying the condition $r_0 < r_1 < 1$ (see (1.7) and (1.8) for the definition and meaning of the number r_0). Therefore, by standard theory of small-norm Riemann–Hilbert problems [16, 17], there exists n_* such that for all $n > n_*$ the R -RH problem is solvable and

$$R(z) = I + R_1(z) + R_2(z) + R_3(z) + \cdots, \quad z \in \mathbb{C} \setminus \Gamma_R, \quad n \geq n_*,$$

where each R_k is of order $O(e^{-kcn})$ and they can be found recursively from

$$R_k(z) = \frac{1}{2\pi i} \int_{\Gamma_R} \frac{[R_{k-1}(\mu)]_-(J_R(\mu) - I)}{\mu - z} d\mu, \quad z \in \mathbb{C} \setminus \Gamma_R, \quad k \geq 1. \quad (4.10)$$

Note that this recurrence also means that

$$R_{k+1}(n) = o(R_k(n)), \quad n \rightarrow \infty.$$

In particular,

$$\begin{aligned} R_1(z; n) &= \frac{1}{2\pi i} \int_{\Gamma_R} \frac{J_R(\mu; n) - I}{\mu - z} d\mu \\ &= \begin{pmatrix} 0 & R_{1,12}(z; n) & 0 & R_{1,14}(z; n) \\ R_{1,21}(z; n) & 0 & R_{1,23}(z; n) & 0 \\ 0 & R_{1,32}(z; n) & 0 & R_{1,34}(z; n) \\ R_{1,41}(z; n) & 0 & R_{1,43}(z; n) & 0 \end{pmatrix}, \end{aligned} \quad (4.11)$$

where

$$\begin{aligned} R_{1,jk}(z; n) &= \frac{1}{2\pi i} \int_{\Gamma'_i} \frac{\mu^n g_{jk}(\mu)}{\mu - z} d\mu, \quad jk = 12, 14, 23, 43, \\ R_{1,jk}(z; n) &= \frac{1}{2\pi i} \int_{\Gamma'_o} \frac{\mu^{-n} g_{jk}(\mu)}{\mu - z} d\mu, \quad jk = 21, 32, 34, 41. \end{aligned} \quad (4.12)$$

Also from (4.10) we can write an expression for R_2 :

$$\begin{aligned} R_2(z; n) &= \frac{1}{2\pi i} \int_{\Gamma_R} \frac{[R_1(\mu; n)]_-(J_R(\mu; n) - I)}{\mu - z} d\mu \\ &= \begin{pmatrix} R_{2,11}(z; n) & 0 & R_{2,13}(z; n) & 0 \\ 0 & R_{2,22}(z; n) & 0 & R_{2,24}(z; n) \\ R_{2,31}(z; n) & 0 & R_{2,33}(z; n) & 0 \\ 0 & R_{2,42}(z; n) & 0 & R_{2,44}(z; n) \end{pmatrix}, \end{aligned}$$

where

$$R_{2,kj}(z; n) = \begin{cases} \sum_{\ell \in \{2,4\}} \frac{1}{2\pi i} \int_{\Gamma'_o} \frac{\mu^{-n} [R_{1,k\ell}(\mu; n)]_- g_{\ell j}(\mu)}{\mu - z} d\mu, & j = 1, k = 1, 3, \\ \sum_{\ell \in \{2,4\}} \frac{1}{2\pi i} \int_{\Gamma'_i} \frac{\mu^n [R_{1,k\ell}(\mu; n)]_- g_{\ell j}(\mu)}{\mu - z} d\mu, & j = 3, k = 1, 3, \\ \frac{1}{2\pi i} \int_{\Gamma'_i} \frac{\mu^n [R_{1,k1}(\mu; n)]_- g_{1j}(\mu)}{\mu - z} d\mu \\ + \frac{1}{2\pi i} \int_{\Gamma'_o} \frac{\mu^{-n} [R_{1,k3}(\mu; n)]_- g_{3j}(\mu)}{\mu - z} d\mu, & k, j = 2, 4. \end{cases}$$

Moreover, using (4.10) and a straightforward calculation one can justify that the matrix structure (i.e., the location of zero and nonzero elements) of R_{2k+1} and R_{2k} , $k \geq 1$, are similar to that of R_1 and R_2 , respectively. It is also straightforward to show that

$$R_{k,ij}(z; n) = \frac{\mathcal{O}(e^{-kcn})}{|z| + 1}, \quad n \rightarrow \infty, \quad k \geq 1,$$

uniformly for $z \in \mathbb{C} \setminus \Gamma_R$, and the positive constant c is the same as in (4.9).

4.2 Asymptotics of h_n

The analysis of the previous section shows that the \mathcal{X} -RH problem has unique solution for all $n > n_*$. We now proceed to the reconstruction of the corresponding functions $\mathcal{X}(z; n)$, $\mathcal{Y}(z; n)$, and to the asymptotics of h_n . To this end we need the asymptotic information about the matrix $P(n)$ which is needed to determine constants C_j that participate in equation (2.49).

Tracing back the Riemann–Hilbert transformations, we find that for $z \in \Omega_0$ we have

$$\mathcal{X}(z; n) = R(z; n) \Lambda(z) \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^{-n} & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^n \end{pmatrix}. \quad (4.13)$$

From (2.48) and (4.13) we conclude that

$$P(n) = R(0; n) \Lambda(0). \quad (4.14)$$

From (4.11) we arrive at the estimate

$$\begin{aligned} R(0; n) &= I + R_1(0; n) + O(e^{-2cn}) \\ &= \begin{pmatrix} 1 & R_{1,12}(0; n) & 0 & R_{1,14}(0; n) \\ R_{1,21}(0; n) & 1 & R_{1,23}(0; n) & 0 \\ 0 & R_{1,32}(0; n) & 1 & R_{1,34}(0; n) \\ R_{1,41}(0; n) & 0 & R_{1,43}(0; n) & 1 \end{pmatrix} + O(e^{-2cn}). \end{aligned} \quad (4.15)$$

Simultaneously, from (4.6) and (4.7) we have

$$\Lambda(0) = \begin{pmatrix} 0 & 0 & 0 & -\alpha(0) \\ -1 & 0 & 0 & 0 \\ 0 & -\frac{1}{\alpha(0)} & 0 & 0 \\ -C_\rho(0)\alpha(0) & 0 & 1 & 0 \end{pmatrix}. \quad (4.16)$$

Equations (4.14), (4.15) and (4.16) yield the following asymptotic formula for $P(n)$,

$$\begin{aligned} P(n) &= \begin{pmatrix} -C_\rho(0)\alpha(0)R_{1,14}(0; n) - R_{1,12}(0; n) & 0 & R_{1,14}(0; n) & -\alpha(0) \\ -1 & -\frac{R_{1,23}(0; n)}{\alpha(0)} & 0 & -\alpha(0)R_{1,21}(0; n) \\ -C_\rho(0)\alpha(0)R_{1,34}(0; n) - R_{1,32}(0; n) & -\frac{1}{\alpha(0)} & R_{1,34}(0; n) & 0 \\ -C_\rho(0)\alpha(0) & -\frac{R_{1,43}(0; n)}{\alpha(0)} & 1 & -\alpha(0)R_{1,41}(0; n) \end{pmatrix} \\ &\quad + O(e^{-2cn}), \end{aligned} \quad (4.17)$$

as $n \rightarrow \infty$.

It is time now for the conditions of Lemma 2.7. We are not going to study each and every condition of this lemma, rather as a case study we consider in particular the condition (2.36):

$$(1 - P_{21}(n))P_{42}(n) + P_{22}(n)P_{41}(n) \neq 0. \quad (4.18)$$

From (4.17) we have that

$$(1 - P_{21}(n))P_{42}(n) + P_{22}(n)P_{41}(n) = -\mathcal{E}(n) + O(e^{-2cn}), \quad n \rightarrow \infty, \quad (4.19)$$

where (cf. (1.6))

$$\mathcal{E}(n) := \frac{2}{\alpha(0)} R_{1,43}(0; n) - C_\rho(0) R_{1,23}(0; n). \quad (4.20)$$

This is when we arrive at condition (1.9) of Theorem 1.1. Indeed, let us suppose that there exist such $C > 0$ and $n_0 \geq n_*$ that

$$|\mathcal{E}(n)| \geq Cr^n, \quad \text{for some } r: r_0 \leq r < 1, \quad \text{and } n > n_0,$$

and choose r_1 in the definition of the constant c so that $r < r_1 < 1$ and $r_1^2 < r$. Then, estimate (4.19) can be rewritten as

$$(1 - P_{21}(n))P_{42}(n) + P_{22}(n)P_{41}(n) = -\mathcal{E}(n)(1 + O(e^{-c_1 n})), \quad n \rightarrow \infty, \quad (4.21)$$

where $c_1 = -\log\left(\frac{r_1^2}{r}\right) > 0$. The last estimate in turn means that there exists such $n_1 \geq n_0$ that for all $n > n_1 + 1$ we should have that condition (4.18) holds for n and $n - 1$ and hence

by Lemma 2.9 and Corollary (2.10), we can uniquely reconstruct the solution of the \mathcal{Y} -RHP, having already the unique solution of the \mathcal{X} -RHP, and, moreover, we could use equation (cf. (2.51)),

$$-\frac{1}{h_{n-1}} = \lim_{z \rightarrow 0} z^{n-1} \mathcal{Y}_{21}(z^{-1}; n). \quad (4.22)$$

for evaluation of the large n behavior of h_n . Let us denote

$$\mathcal{A}(z; n) := P^{-1}(n) \mathcal{X}(z; n) \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^n & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^{-n} \end{pmatrix}, \quad (4.23)$$

and also let us define the matrix $\mathcal{B}(n)$ in the following expansion for $\mathcal{A}(z; n)$, which is equivalent to **RH- \mathcal{X} 4**:

$$\mathcal{A}(z; n) = I + \mathcal{B}(n)z + \mathcal{O}(z^2), \quad z \rightarrow 0. \quad (4.24)$$

Therefore by (2.29), (2.30), (2.31) and (4.23) we can write

$$\overset{\circ}{\mathcal{X}}(z, n) = \begin{pmatrix} C_1(n) & 1 & C_3(n) & 0 \\ C_2(n) & 0 & C_4(n) & 1 \end{pmatrix} \mathcal{A}(z; n) \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & z^{-n} & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & z^n \end{pmatrix}. \quad (4.25)$$

Using (2.19) and (4.25) we can write

$$\mathcal{Y}_{21}(z^{-1}; n) = \overset{\circ}{\mathcal{X}}_{22}(z; n) = C_2(n) \mathcal{A}_{12}(z; n) z^{-n} + C_4(n) \mathcal{A}_{32}(z; n) z^{-n} + \mathcal{A}_{42}(z; n) z^{-n}. \quad (4.26)$$

From (4.24) we have

$$z^{-n} \mathcal{A}(z; n) = z^{-n} I + z^{-n+1} \mathcal{B}(n) + \mathcal{O}(z^{-n+2}), \quad z \rightarrow 0.$$

Therefore, as $z \rightarrow 0$

$$z^{-n} \mathcal{A}_{ij}(z; n) = \begin{cases} z^{-n+1} \mathcal{B}_{ij}(n) + \mathcal{O}(z^{-n+2}), & i \neq j, \\ z^{-n} + z^{-n+1} \mathcal{B}_{ii}(n) + \mathcal{O}(z^{-n+2}), & i = j. \end{cases} \quad (4.27)$$

Therefore by (4.22), (4.26) and (4.27) we have

$$-\frac{1}{h_{n-1}} = C_2(n) \mathcal{B}_{12}(n) + C_4(n) \mathcal{B}_{32}(n) + \mathcal{B}_{42}(n). \quad (4.28)$$

Let us denote the coefficients in the expansions of $R(z; n)$ and $\Lambda(z)$, as $z \rightarrow 0$, by

$$\begin{aligned} R(z; n) &= R(0; n) + R^{(1)}(n)z + R^{(2)}(n)z^2 + \mathcal{O}(z^3), \\ \Lambda(z) &= \Lambda(0) + \Lambda^{(1)}z + \Lambda^{(2)}z^2 + \mathcal{O}(z^3). \end{aligned}$$

From (4.24), (4.13), and (4.14) we have

$$\mathcal{B}(n) = \Lambda^{-1}(0) R^{-1}(0; n) R^{(1)}(n) \Lambda(0) + \Lambda^{-1}(0) \Lambda^{(1)}. \quad (4.29)$$

Note that

$$R^{(1)}(n) = \frac{1}{2\pi i} \int_{\Gamma_R} (J_R(\mu; n) - I) \frac{d\mu}{\mu^2} + \mathcal{O}(e^{-2cn}), \quad R^{-1}(0; n) = I - R_1(0; n) + \mathcal{O}(e^{-2cn}),$$

as $n \rightarrow \infty$. More precisely, we have

$$R^{(1)}(n) = \begin{pmatrix} 0 & R_{12}^{(1)}(n) & 0 & R_{14}^{(1)}(n) \\ R_{21}^{(1)}(n) & 0 & R_{23}^{(1)}(n) & 0 \\ 0 & R_{32}^{(1)}(n) & 0 & R_{34}^{(1)}(n) \\ R_{41}^{(1)}(n) & 0 & R_{43}^{(1)}(n) & 0 \end{pmatrix} + \mathcal{O}(e^{-2cn}), \quad (4.30)$$

as $n \rightarrow \infty$, where

$$\begin{aligned} R_{jk}^{(1)}(n) &= \frac{1}{2\pi i} \int_{\Gamma'_i} \mu^{n-2} g_{jk}(\mu) d\mu, & jk &= 12, 14, 23, 43, \\ R_{jk}^{(1)}(n) &= \frac{1}{2\pi i} \int_{\Gamma'_o} \mu^{-n-2} g_{jk}(\mu) d\mu, & jk &= 21, 32, 34, 41, \end{aligned} \quad (4.31)$$

and

$$R^{-1}(0; n) = \begin{pmatrix} 1 & -R_{1,12}(0; n) & 0 & -R_{1,14}(0; n) \\ -R_{1,21}(0; n) & 1 & -R_{1,23}(0; n) & 0 \\ 0 & -R_{1,32}(0; n) & 1 & -R_{1,34}(0; n) \\ -R_{1,41}(0; n) & 0 & -R_{1,43}(0; n) & 1 \end{pmatrix} + \mathcal{O}(e^{-2cn}), \quad (4.32)$$

$$n \rightarrow \infty.$$

Regarding the coefficients of $\Lambda(z)$, we will actually only need the zero-term, $\Lambda(0)$, which has already been presented in (4.16).

From (4.29), (4.30), (4.32) and (4.16) we find that

$$\begin{aligned} \mathcal{B}_{12}(n) &= \frac{R_{23}^{(1)}(n)}{\alpha(0)} + \mathcal{O}(e^{-2cn}), & \mathcal{B}_{32}(n) &= C_\rho(0)R_{23}^{(1)}(n) - \frac{R_{43}^{(1)}(n)}{\alpha(0)} + \mathcal{O}(e^{-2cn}), \\ \mathcal{B}_{42}(n) &= -\frac{1}{\alpha^2(0)}(R_{1,12}(0; n)R_{23}^{(1)}(n) + R_{1,14}(0; n)R_{43}^{(1)}(n)) + \mathcal{O}(e^{-3cn}). \end{aligned} \quad (4.33)$$

Note that $\mathcal{B}_{12}(n)$, $\mathcal{B}_{32}(n)$ are of order $\mathcal{O}(e^{-cn})$, while $\mathcal{B}_{42}(n)$ is of order $\mathcal{O}(e^{-2cn})$.

Revisiting (2.31) we have

$$\begin{pmatrix} 1 & C_1(n) & 0 & C_3(n) \\ 0 & C_2(n) & 1 & C_4(n) \end{pmatrix} P(n) = \begin{pmatrix} C_1(n) & 1 & C_3(n) & 0 \\ C_2(n) & 0 & C_4(n) & 1 \end{pmatrix}.$$

In particular, in view of (4.18), we can write the following two equations for the constants C_2 and C_4

$$\begin{aligned} C_2(n)P_{21}(n) + P_{31}(n) + C_4(n)P_{41}(n) &= C_2(n), \\ C_2(n)P_{22}(n) + P_{32}(n) + C_4(n)P_{42}(n) &= 0. \end{aligned}$$

Solving for C_2 and C_4 we find

$$\begin{aligned} C_2(n) &= \frac{P_{42}(n)P_{31}(n) - P_{41}(n)P_{32}(n)}{(1 - P_{21}(n))P_{42}(n) + P_{41}(n)P_{22}(n)}, \\ C_4(n) &= -\frac{P_{22}(n)P_{31}(n) + [1 - P_{21}(n)]P_{32}(n)}{(1 - P_{21}(n))P_{42}(n) + P_{41}(n)P_{22}(n)}. \end{aligned}$$

From (4.17) and (4.21) we have

$$C_2(n) = \frac{C_\rho(0)}{\mathcal{E}(n)}(1 + \mathcal{O}(e^{-c_1 n})), \quad (4.34)$$

and

$$C_4(n) = -\frac{2}{\alpha(0)\mathcal{E}(n)}(1 + \mathcal{O}(e^{-c_1 n})). \quad (4.35)$$

Combining (4.28), (4.33), (4.34) and (4.35) we obtain

$$h_{n-1} = -\alpha(0) \frac{\mathcal{E}(n)}{\frac{2}{\alpha(0)}R_{43}^{(1)}(n) - C_\rho(0)R_{23}^{(1)}(n)}(1 + \mathcal{O}(e^{-c_1 n})), \quad n \rightarrow \infty. \quad (4.36)$$

Note that from (4.12) and (4.31) we have

$$\begin{aligned} R_{1,jk}(0; n) &= R_{jk}^{(1)}(n+1), & \text{for } jk = 12, 14, 23, 43, \\ R_{1,jk}(0; n) &= R_{jk}^{(1)}(n-1), & \text{for } jk = 21, 32, 34, 41. \end{aligned}$$

This allows us to rewrite (4.36) as

$$h_{n-1} = -\alpha(0) \frac{\mathcal{E}(n)}{\mathcal{E}(n-1)}(1 + \mathcal{O}(e^{-c_1 n})), \quad n \rightarrow \infty. \quad (4.37)$$

This concludes the proof of Theorem 1.1.

Let \mathfrak{C} denote the class of symbol pairs $(\phi, d\phi)$, where ϕ and d satisfy the properties mentioned in Theorem 1.1, for which the corresponding matrix $P(n)$ satisfies the condition (4.18) for sufficiently large n . Now, consider the subclass $\mathfrak{C}_0 \subset \mathfrak{C}$, which includes symbol pairs $(\phi, d\phi)$ for which $\mathcal{E}(n)$ satisfies inequality (1.9), or rather its generalization (1.11) for sufficiently large n (see also Section 5.5). For the purposes of this paper, it is worthwhile to prove that the subclass \mathfrak{C}_0 is not empty by providing an explicit example. To this end, let us consider the symbol pairs $(\phi, d\phi)$, where

$$\phi(z) = \left(\frac{z-b}{z-a} \right)^\alpha, \quad (4.38)$$

and

$$d(z) = \left(\frac{z-b_1}{z-a_1} \right)^{\alpha_1} \left(\frac{a_1 z - 1}{b_1 z - 1} \right)^{\alpha_1}. \quad (4.39)$$

We assume that $0 < a < b < a_1 < b_1 < 1$, and $|\Re \alpha_1| < 1$. Then, for $\mathcal{E}(n)$ defined in (4.20), by standard Watson lemma type asymptotic analysis of integrals we will arrive at

$$\mathcal{E}(n) = \kappa b_1^{n-\alpha_1} n^{\alpha_1-1} \left(1 + \mathcal{O}\left(\frac{1}{n}\right) \right), \quad n \rightarrow \infty, \quad (4.40)$$

where

$$\begin{aligned} \kappa &= -\frac{i}{\pi} e^{-i\pi\alpha_1} \Gamma(1-\alpha_1) (b_1 - a_1)^{\alpha_1} \left(\frac{1-bb_1}{1-ab_1} \right)^\alpha \left(\frac{a_1}{b_1} \right)^{-\alpha_1} \\ &\quad \times \left[e^{-i\pi\alpha_1} \frac{i}{\pi} \int_{1/b_1}^{1/a_1} \left(\frac{z-b_1}{z-a_1} \right)^{\alpha_1} \left(\frac{1/a_1 - z}{z - 1/b_1} \right)^{-\alpha_1} \frac{z+b_1}{z-b_1} \frac{dz}{z} - 1 \right], \end{aligned} \quad (4.41)$$

and

$$\arg \left(\frac{z - b_1}{z - a_1} \right) = \arg \left(\frac{1/a_1 - z}{z - 1/b_1} \right) = 0, \quad z \in [1/b_1, 1/a_1].$$

We argue, that for generic values of a, b, a_1, b_1, α and α_1 the number $\kappa \neq 0$. Indeed, let us take, for instance, α_1 real, i.e.,

$$-1 < \alpha_1 < 1.$$

Then, the integrand in (4.41) is strictly positive and hence, except for $\alpha_1 = \pm 1/2$, the first term in the brackets of (4.41) has nonzero imaginary part and thus can not cancel the second term. It is only needed now to notice that in the case under consideration $r_0 = b_1$ and hence inequality $\kappa \neq 0$ implies that $\mathcal{E}(n)$ corresponding to the weights (4.38) and (4.39) satisfies the condition (1.11) indicated in Remark 1.2 with $\varkappa(n) = n^{\alpha_1-1}$.

We conclude the discussion of this example by indicating the asymptotic behavior of the corresponding norm parameter h_{n-1} and the determinant D_n . To this end we notice that, in fact, the estimate (4.40) can be extended to the full asymptotic series,

$$\mathcal{E}(n) \sim \kappa b_1^{n-\alpha_1} n^{\alpha_1-1} \left(1 + \sum_{k=1}^{\infty} d_k n^{-k} \right), \quad n \rightarrow \infty.$$

Therefore, from (4.37) we should have that

$$h_{n-1} = -b_1 \left(\frac{n}{n-1} \right)^{\alpha_1-1} \left(1 + O\left(\frac{1}{n^2} \right) \right), \quad n \rightarrow \infty,$$

and hence

$$D_n = C(-b_1)^n n^{\alpha_1-1} \left(1 + O\left(\frac{1}{n} \right) \right), \quad n \rightarrow \infty,$$

for some constant C which is yet to be determined. It is interesting that this formula indicates an oscillatory behavior of the determinant D_n for large n . This fact is confirmed numerically.

Remark 4.1. From (2.5), (2.19), (2.29), and (2.30) we have the following representation for the orthogonal polynomials $\mathcal{P}_n(z)$ in terms of the solution $\mathcal{X}(z; n)$ of the \mathcal{X} -RHP,

$$\mathcal{P}_n(z) = \mathcal{X}_{11}(z; n) + C_1(n)\mathcal{X}_{21}(z; n) + C_3(n)\mathcal{X}_{41}(z; n).$$

The asymptotic results concerning the function $\mathcal{X}(z; n)$ obtained in this section can be translated to the large n asymptotic formulae for the polynomials $\mathcal{P}_n(z)$. Indeed, skipping the rather tedious though straightforward calculations, we arrive at the following asymptotics for $\mathcal{P}_n(z)$ on the unit circle:

$$\begin{aligned} \mathcal{P}_n(z) &= \frac{\beta_+(z)(2\alpha(0)C_{\rho,+}(z) - C_{\rho}(0)\alpha(0))}{\mathcal{E}(n)} \\ &\quad + \alpha_-(z)z^n \left(1 + \frac{C_{\rho}(0)\alpha(0)R_{1,21}(z; n) - 2R_{1,41}(z; n)}{\mathcal{E}(n)} \right) + \mathcal{O}(e^{-c_1 n}), \end{aligned} \quad (4.42)$$

as $n \rightarrow \infty$. While in the interior and exterior of the unit circle we have the following asymptotic formulae for $\mathcal{P}_n(z)$:

$$\begin{aligned} \mathcal{P}_n(z) &= \frac{\beta(z)(2\alpha(0)C_{\rho}(z) - C_{\rho}(0)\alpha(0))}{\mathcal{E}(n)} + \mathcal{O}(e^{-c_1 n}), \quad |z| < 1, \\ \mathcal{P}_n(z) &= \alpha(z)z^n \left(1 + \frac{C_{\rho}(0)\alpha(0)R_{1,21}(z; n) - 2R_{1,41}(z; n)}{\mathcal{E}(n)} + \mathcal{O}(e^{-2c_1 n}) \right), \quad |z| > 1, \end{aligned}$$

as $n \rightarrow \infty$ (compare with (4.42)).

Remark 4.2. In a similar fashion, we could have obtained the large n asymptotic expression for the norm h_n of orthogonal polynomials (3.1). In that case, the methods of this section would work analogously if the function

$$d := -\phi^{-1}(z)\tilde{u}(z) \equiv \phi^{-1}(z) \int_a^b \frac{t^{s-1}w(t)}{1-tz} dt,$$

satisfies $d\tilde{d} \equiv 1$ on the unit circle, or

$$\phi(z)\tilde{\phi}(z) = \int_a^b \frac{t^{s-1}w(t)}{1-tz} dt \int_a^b \frac{t^{s-1}w(t)}{1-tz^{-1}} dt, \quad z \in \mathbb{T}. \quad (4.43)$$

5 Remaining open questions

We consider this work as a starting point of a long term research project. There are many challenging technical as well as conceptual open questions related to the Riemann–Hilbert formalism we are suggesting in this paper. Here we highlight some of them that we consider the most pressing.

5.1 Derivation of the relevant Christoffel–Darboux formulae and differential identities

Our main objective in this paper has been to develop a 4×4 steepest descent analysis for certain Toeplitz+Hankel determinants and we have achieved that. However, to obtain the asymptotics of $D_n(\phi, d\phi, 1, 1)$ one has to derive suitable differential identities. We propose that the differential identity has to be with respect to the parameters α_i in the function d given by (4.1). Thus, one has to perform m integrations in the parameters α_i , $1 \leq i \leq m$. Note that for $\alpha_1 = \alpha_2 = \dots = \alpha_m = 0$, we have $d \equiv 1$ and hence $\phi = w$. Hence the starting point of the integration in α_1 could be taken from the results of [2].⁵ Integration of the differential identity in α_1 will provide us with an asymptotic expression for $D_n(\phi, d_1\phi, 1, 1)$, which also serves as the starting point of integration in α_2 . Thus we can find asymptotics of $D_n(\phi, d_1d_2\phi, 1, 1)$ which also serves as the starting point of integration in α_3 , and so on. Repeating this procedure will finally lead us to the asymptotics of $D_n(\phi, d\phi, 1, 1)$.

In order to derive the differential identities mentioned above, one has to find recurrence relations and prove a Christoffel–Darboux formula for the polynomials (2.1) and follow a path similar to that introduced by I. Krasovsky in [20]. Note that the recurrence relations can be found by analyzing the function $\mathcal{M}(z; n) := \mathcal{X}(z; n+1)\mathcal{X}^{-1}(z; n)$, which is holomorphic in $\mathbb{C} \setminus \{0\}$ and can be globally determined by its singular parts at zero and infinity.

5.2 Extension of the analysis to general offset values $r, s \neq 1$

If we lift the restriction $r = s = 1$ then in the jump matrix of the \mathcal{X} -RHP the functions $\phi(z)$ and $w(z)$ should be replaced by $z^{1-r}\phi(z)$ and $z^{1-s}w(z)$, respectively. This, in turn would raise a serious question about solvability of the Λ -RHP. Indeed, for instance, we would not be able to define the function α in (4.3) and hence to factorize the jump matrix (4.2). The way out of this difficulty could be the use of the relation between the determinants $D_n(\phi, w; r, s)$ with different values of r, s or both. Such relations are well known in the pure Toeplitz case (for example see [13, Lemma 2.4]). However, for general Toeplitz+Hankel determinants they are yet to be found.

⁵Although the authors in [2] do not particularly study the asymptotics of $D_n(\phi, \phi, 1, 1)$, this asymptotics is achievable by their methods.

Another way to approach the problem could be to develop the so-called Bäcklund–Schlesinger transformations of the \mathcal{X} -RHP itself. That is, the transformations of the form,

$$\mathcal{X}(z) \mapsto R(z)\mathcal{X}(z),$$

where $R(z)$ is a properly chosen rational function for which the above transformation results in the desired shifting of the parameters r and s . Also, one can try to allow the matrix $P(n)$ in the setting of the \mathcal{X} -RHP to depend on z .

It is worth to point out that the problem with the extension of our scheme to the general values of r and s is not actually a problem of the setting of the relevant Riemann–Hilbert problem. Indeed, it is rather the question of the correct way to approach to its asymptotic solution. Let us demonstrate this point by considering the pure Toeplitz case.

Assume that $D_n(\phi, 0; r, 0) \neq 0$ for all n , so that $\mathcal{P}_n(z)$ and, correspondingly, the solution $\mathcal{Y}(z; n)$ of the \mathcal{Y} -RHP exist for all n . Put,

$$\mathcal{X}(z; n) := P_\infty^{-1}(n) \begin{pmatrix} \mathcal{Y}_{11}(z; n) & 0 & 0 & \tilde{\mathcal{Y}}_{12}(z; n) \\ 0 & \tilde{\mathcal{Y}}_{11}(z; n) & \mathcal{Y}_{12}(z; n) & 0 \\ 0 & \tilde{\mathcal{Y}}_{21}(z; n) & \mathcal{Y}_{22}(z; n) & 0 \\ \mathcal{Y}_{21}(z; n) & 0 & 0 & \tilde{\mathcal{Y}}_{22}(z; n) \end{pmatrix},$$

where the normalizing matrix $P_\infty(n)$ is given by

$$P_\infty(n) = \begin{pmatrix} 1 & 0 & 0 & \mathcal{Y}_{12}(0; n) \\ 0 & \mathcal{Y}_{11}(0; n) & 0 & 0 \\ 0 & \mathcal{Y}_{21}(0; n) & 1 & 0 \\ 0 & 0 & 0 & \mathcal{Y}_{22}(0; n) \end{pmatrix},$$

which is invertible for generic ϕ . It is straightforward to check that the so defined 4×4 matrix-valued function \mathcal{X} solves the \mathcal{X} -Riemann–Hilbert problem, **RH- \mathcal{X} 1–RH- \mathcal{X} 4** with $w \equiv 0$, and $\phi(z)$ replaced by $z^{1-r}\phi(z)$. Take now $r = 0$. That is, let us consider the standard orthogonal polynomials on the circle with the weight $\phi(z)$ having zero winding number. Then, from the standard 2×2 Riemann–Hilbert formalism [1], we know everything about the asymptotic behavior of the corresponding orthogonal polynomials $\mathcal{P}_n(z)$ and hence we know asymptotic solution of the \mathcal{X} -Riemann–Hilbert problem corresponding to $D_n(\phi, 0; 0, 0)$. And, we know this in spite of the fact that the approach developed in the body of this work can not be satisfactorily used for the case $r = s = 0$. We believe that this observation might entail a useful hint on how to modify our Riemann–Hilbert approach to deal with general values of r and s .

5.3 Extension of the Riemann–Hilbert analysis of Section 3 for more general choices of I

We recall that our Riemann–Hilbert analysis of Section 3, with minimal modifications, naturally extends to the following three cases as well: *i*) $-1 < a < b < 0$, *ii*) $-\infty < a < b < -1$, and *iii*) $1 < a < b < \infty$. A natural first step in generalizing the results of Section 3 beyond the above cases is considering the case where I has 0 as an end point. This would slightly affect the analysis as one has to take into account the behavior of w at 0 in the set up of the 2×4 and the subsequent Riemann–Hilbert problems.

The other interesting case to be studied is when I intersects the unit circle. Clearly, in this case one has to perform local analysis in a neighborhood of the intersection point(s) of I and the unit circle. Although these local constructions are reminiscent of what one does near the Fisher–Hartwig singularities or the endpoints of the support of the symbol, here even if the

possible intersection points ± 1 , are regular points for non-FH symbols ϕ and w , one has to still perform local analysis due to collision of the supports of ϕ and w .

Another possible generalization would be to consider I to be the union of two symmetric intervals with respect to the unit circle, i.e., $I = [a, b] \cup [b^{-1}, a^{-1}]$. This generalization should be accessible by slight modification of our approach explained in Section 3. However, generalization to the case where I is a union of two non-symmetrical intervals with respect to the unit circle needs a more special treatment.

5.4 Extension to Fisher–Hartwig symbols

One can study the large- n asymptotics of determinant $D_n(\phi, d\phi, 1, 1)$ (and with increasing effort $D_n(\phi, d\phi, r, s)$ for fixed $r, s \in \mathbb{Z}$) assuming that ϕ possesses Fisher–Hartwig singularities $\{z_i\}_{i=1}^m$ on the unit circle. It is in fact in this level of generality that E. Basor and T. Ehrhardt have been able to compute the asymptotics of $D_n(\phi, d\phi, 0, 1)$, $D_n(-\phi, d\phi, 0, 1)$, $D_n(\phi, dz\phi, 0, 1)$, and $D_n(-z\phi, d\phi, 0, 1)$ via the operator-theoretic methods in [2]. However, the authors in [2] further require that the Fisher–Hartwig part of ϕ be *even*. In fact they used some results of the work [13] to prove their asymptotic formulas for Toeplitz+Hankel determinants, and for this reason they *inherited* the evenness assumption from the work [13] where the authors needed evenness of ϕ in their 2×2 setting to relate Hankel and Toeplitz+Hankel determinants to a Toeplitz determinant with symbol ϕ .

From a Riemann–Hilbert perspective, in the presence of Fisher–Hartwig singularities, one has to construct the 4×4 local parametrices near the points z_i . Expectedly, these local parametrices must be expressed in terms of confluent hypergeometric functions as suggested by [13]. We have not yet worked out the details of this construction but we believe that it should be well within reach. It would be methodologically important to achieve the results obtained from operator-theoretic tools via the Riemann–Hilbert approach as well. Moreover, we expect that the evenness of the Fisher–Hartwig part of ϕ would not play a role in our 4×4 setting, and in that sense there are reasonable prospects of generalizing the results of [2] to symbols ϕ with non-even Fisher–Hartwig part.

5.5 Characterization of generic classes of Szegő-type symbol pairs $(\phi, d\phi)$, with $d\tilde{d} \equiv 1$ on the unit circle

Take one of the six conditions of Lemma 2.7. Denote by $\mathfrak{P}(n)$ and $\mathcal{E}(n)$, respectively, the corresponding nonzero quantity and its leading order term in the large n asymptotic expansion. Consider the class of symbol pairs $(\phi, d\phi)$, where ϕ and d satisfy the properties mentioned in Theorem 1.1. Within this class, take the subclass \mathfrak{C} of symbol pairs for which the elements of the corresponding matrix $P(n)$ satisfy $\mathfrak{P}(n) \neq 0$ for sufficiently large n . Also consider the subclass $\mathfrak{C}_0 \subseteq \mathfrak{C}$ of symbols pairs for which $\mathcal{E}(n)$ satisfies (1.11) for sufficiently large n . One should be able to find the asymptotics of the norm h_n of orthogonal polynomials in a similar fashion as presented in Section 4.2 in terms of $\mathcal{E}(n)$, assuming that $(\phi, d\phi) \in \mathfrak{C}_0$. It would be very interesting to completely or partially characterize the classes of symbol pairs \mathfrak{C} and \mathfrak{C}_0 , corresponding to each one of the six conditions of Lemma 2.7.

Moreover, although we have provided explicit examples for a class of Szegő-type functions d which further satisfy $d\tilde{d} \equiv 1$ on the unit circle (see the beginning of Section 4), a complete characterization of such functions is currently unknown to the authors. Also, regarding what we discussed in Section 3, for a given Szegő-type symbol ϕ , we are very interested to find the associated class of weights w , supported on the interval, for which the equality (4.43) holds.

5.6 Characteristic polynomial of a Hankel matrix

As mentioned in the Introduction, arguably the most important motivation behind studying the asymptotics of Toeplitz+Hankel determinants is to study the large n asymptotics of the eigenvalues of the matrix $H_n[w]$. We recall that the characteristic polynomial $\det(H_n[w] - \lambda I)$ is indeed the Toeplitz+Hankel determinant $D_n(-\lambda, w, 0, 0)$. In this case the associated Λ -model Riemann–Hilbert problem needs a special treatment. In a sense it is a simpler problem as the symbol ϕ is identically equal to a constant, but more complicated – compared to the situation in Section 4 – as it does not enjoy $J_{\Lambda,23}(z) = 0$. In any case, the solution to this model problem provides us with the constant term in the asymptotics of $D_n(-\lambda, w, 0, 0)$, and in the case of Fisher–Hartwig weight w , one can hope to obtain the leading terms of this asymptotic expansion (up to the constant term, viz. the solution of the Λ -model problem) from the local analysis near the Fisher–Hartwig singularities.⁶ This last point is yet another motivation to pursue the goals of Section 5.4.

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⁶The very recent work [18] shows that indeed one can explicitly describe the leading terms of the asymptotics of $D_n(-\lambda, w, 0, 0)$, in the case of w having jump discontinuities and λ sufficiently large.

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