A pair of commuting hypergeometric operators on the complex plane and bispectrality

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We consider the standard hypergeometric differential operator \( D \) regarded as an operator on the complex plane \( \mathbb{C} \) and the complex conjugate operator \( \overline{D} \). These operators formally commute and are formally adjoint one to another with respect to an appropriate weight. We find conditions when they commute in the Nelson sense and write explicitly their joint spectral decomposition. It is determined by a two-dimensional counterpart of the Jacobi transform (synonyms: generalized Mehler–Fock transform, Olevskii transform). We also show that the inverse transform is an operator of spectral decomposition for a pair of commuting difference operators defined in terms of shifts in imaginary direction.

1. Introduction

1.1. Spectral problem. Denote by \( \hat{\mathbb{C}} \) the complex plane without the points 0 and 1, by \( \mathcal{D}(\hat{\mathbb{C}}) \) the space of smooth compactly supported functions on \( \hat{\mathbb{C}} \). Denote by \( d\overline{z} \) the standard Lebesgue measure on \( \mathbb{C} \).

Fix real \( a \) and \( b \). Consider the following measure on \( \hat{\mathbb{C}} \):

\[
\mu_{a,b}(z) \ d\overline{z} := |z|^{2a+2b-2}|1-z|^{2a-2b} \ d\overline{z}
\]

and the corresponding space \( L^2(\mathbb{C}, \mu_{a,b}) \),

\[
\langle f, g \rangle = \int_{\hat{\mathbb{C}}} f(z) \overline{g(z)} \mu_{a,b}(z) \ d\overline{z}.
\]

Consider the following pair of differential operators in the space \( L^2(\mathbb{C}, \mu_{a,b}) \):

\[
\mathcal{D} := z(1-z) \frac{\partial^2}{\partial z^2} + (a+b-(2a+1)z) \frac{\partial}{\partial z} - a^2;
\]

\[
\overline{\mathcal{D}} := \overline{z}(1-\overline{z}) \frac{\partial^2}{\partial \overline{z}^2} + (a+b-(2a+1)\overline{z}) \frac{\partial}{\partial \overline{z}} - a^2.
\]

These operators formally commute, i.e.,

\[
\mathcal{D} \overline{\mathcal{D}} f = \overline{\mathcal{D}} \mathcal{D} f, \quad \text{where } f \in \mathcal{D}(\hat{\mathbb{C}}).
\]

A straightforward calculation shows that they are formally adjoint,

\[
\langle \mathcal{D} f, g \rangle = \langle f, \overline{\mathcal{D}} g \rangle, \quad \text{where } f, g \in \mathcal{D}(\hat{\mathbb{C}}).
\]

Therefore the operators \( \frac{1}{i}(\mathcal{D} + \overline{\mathcal{D}}) \), \( \frac{1}{i}(\mathcal{D} - \overline{\mathcal{D}}) \) are symmetric on the domain \( \mathcal{D}(\hat{\mathbb{C}}) \).

The purpose of this paper is to construct an explicit spectral decomposition of this pair, i.e., a unitary operator \( U \), which diagonalizes both operators \( \mathcal{D}, \overline{\mathcal{D}} \).

As we know after the famous work of Edward Nelson [33], 1959, (see, also [42], Sect. VIII.5) a question about commutativity of two unbounded self-adjoint operators can be highly nontrivial. Recall that two self-adjoint operators \( A, B \) commute
Figure 1. To Theorem 1.1. The domain $\Pi$ of commutativity, and the domain $\Pi_{\text{cont}} \subset \Pi$, where the spectrum is purely continuous.

if they can be simultaneously realized as operators of multiplication by functions in some $L^2$. Equivalently, the corresponding one-parametric groups commute:

$$e^{isA}e^{itB} = e^{itB}e^{isA},$$

where $s, t$ in $\mathbb{R}$.

Equivalently, resolvents $(A - \lambda)^{-1}$ and $(B - \mu)^{-1}$, commute. However these properties do not follow from the identity $AB = BA$ and are difficult for a verification. There are some useful sufficient conditions and necessary conditions for commutativity (for necessity we use the result of Kostyuchenko and Mityagin [23-24]), but quite often a question remains to be heavy.

Define two domains $\Pi \supset \Pi_{\text{cont}}$ of the parameters $(a, b)$:

(1.4) $\Pi : 0 < a + b < 2, \quad -1 < a - b < 1.$

(1.5) $\Pi_{\text{cont}} : 0 \leq a \leq 1, \quad 0 \leq b \leq 1, \quad \text{and} \quad (a, b) \neq (\pm 1, \pm 1), (\pm 1, \mp 1).$

Theorem 1.1. The operators $\frac{1}{2}(D + \bar{D}), \frac{1}{2i}(D - \bar{D})$ admit extensions to a pair of commuting self-adjoint operators if and only if $(a, b) \in \Pi$.

Next, we define a natural domain for our operators. Consider the subspace $\mathcal{R}_{a,b}(\hat{C}) \subset L^2(\hat{C}, \mu_{a,b})$ consisting of smooth functions $f$ on $\hat{C}$ satisfying the following conditions:

1°. In a neighborhood of $z = 0$ a function $f$ has an expansion of the form:

(1.6) $$f(z) = \begin{cases} 
\alpha(z) + \beta(z) |z|^{2-2a-2b}, & \text{if } a + b \neq 1; \\
\alpha(z) + \beta(z) \ln |z|, & \text{if } a + b = 1,
\end{cases}$$

where $\alpha(z), \beta(z)$ are smooth functions.

2°. In a neighborhood of $z = 1$ a function $f$ has an expansion of the form:

(1.7) $$f(z) = \begin{cases} 
\gamma(z) + \delta(z) |z - 1|^{2b-2a}, & \text{if } a - b \neq 0; \\
\gamma(z) + \delta(z) \ln |z - 1|, & \text{if } a - b = 0,
\end{cases}$$

where $\gamma(z), \delta(z)$ are smooth.

4A famous example is a problem, see [12], which was raised by Irving Segal in 1958 and which was discussed during almost 30 years: Let $\Omega$ be an open connected domain in $\mathbb{R}^n$. Assume that the operators $i\partial/\partial x_k$ in $\mathcal{D}(\Omega)$ admit commuting self-adjoint extensions. Is it correct that $\Omega$ is essentially a fundamental domain of $\mathbb{R}^n$ with respect to a certain discrete group? The answer is affirmative.

5If $(a, b) \notin \Pi$, then $\mathcal{R}_{a,b}(\hat{C})$ is not contained in $L^2(\mathcal{C}, \mu_{a,b})$.

6Boundary conditions in this spirit sometimes arise in spectral theory of ordinary differential operators $D$ for operators with deficiency indices $(1, 1)$ or $(2, 2)$, see, e.g., [36], Section 1.
For each \( p, q, N \) we have
\[
\frac{\partial^p q f}{\partial z^p} = O \left( |z|^{-2a-p-q} (\ln |z|)^{-N} \right) \quad \text{as } z \to \infty.
\]

**Theorem 1.2.**

a) For \((a, b) \in \Pi \) the operators \( \frac{1}{2}(D + \overline{D}) \) are essentially self-adjoint on \( \mathcal{D} \) and commute in the Nelson sense.

b) If \((a, b) \in \Pi_{\text{cont}} \), then the spectrum of the problem
\[
Df = \zeta f, \quad \overline{D}f = \overline{\zeta} f
\]
is multiplicity free and consists of \( \zeta \) having the form
\[
\zeta = \left( \frac{k+i\ell}{2} \right)^2, \quad \text{where } k \in \mathbb{Z}, \ s \in \mathbb{R}.
\]
If \((a, b) \in \Pi \setminus \Pi_{\text{cont}} \), then the spectrum consists on the same set plus one eigenvalue \( \zeta_0 > 0 \).

Let us explain the obstacle for commutativity. Consider a second order differential operator \( D \) on an interval. For each \( \zeta \in \mathbb{R} \) the differential equation \( Df = \zeta f \) has two solutions, and we can select generalized eigenfunctions of \( D \) as solutions that have \( L^2 \)- or almost \( L^2 \)-asymptotics at the ends of the interval. In our case the system \( (1.9) \) locally has 4 solutions. Furthermore, \( \mathbb{C} \) is not simply connected, solutions are ramified at 0, 1, \( \infty \). As a result there are few single valued solutions and we have no freedom for selection of asymptotics. Such considerations (see Section 4) allow to establish necessity of the conditions of Theorem 1.1.

Unfortunately, we do not know an a priori proof of sufficiency and obtain it as a byproduct of the explicit joint spectral decomposition of the operators \( D, \overline{D} \). Such detour makes our work long and requires numerous explicit calculations and estimates.

### 1.2. The index hypergeometric transform.

Our work is a counterpart of the following classical topic. Consider the hypergeometric differential operator
\[
D := x(x+1) \frac{d^2}{dx^2} + \left( (a+b) + (2a+1)x \right) \frac{d}{dx} + a^2
\]
on the half-line \( \mathbb{R}_+ \), i.e., \( x > 0 \). Consider the integral operator
\[
I_{a,b}f(s) := \frac{1}{\Gamma(a+b)} \int_0^\infty f(x) \binom{x}{a-i\ell, a+i\ell; x} \left[ x^{a+b-1}(1+x)^{a-b} \right] dx.
\]
Then \( I_{a,b} \) is a unitary operator
\[
L^2 \left[ \mathbb{R}_+, x^{a+b+1}(1+x)^{a-b} \right] \to L^2 \left[ \mathbb{R}_+, \pi^{-1} \left| \frac{\Gamma(a+i\ell)\Gamma(b+i\ell)}{\Gamma(2i\ell)} \right|^2 ds \right].
\]
The operator \( I_{a,b} \) sends \( D \) to the multiplication by \( s^2 \), see [15, 33, 40, 22, 21, 34, 37]. This transform is known as ‘the generalized Mehler-Fock transform’, ‘the Olevskii transform’, or ‘the Jacobi transform’.

Such operators arise in a natural way in the analysis on rank one Riemannian symmetric spaces, on the other hand they are special cases of multi-dimensional

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\( ^7 \) A special case \( a = 1/2, b = 1 \) of this transform was discovered by Gustav Mehler in 1881, the general transform was obtained by Weyl [45] in 1910. The \( I_{a,b} \) is a representative of a large family of index integral transforms, which involve indices of hypergeometric functions, see numerous examples in [10, 19, 39].
Harish-Chandra spherical transforms and more general Heckman–Opdam integral transforms, which arise as spectral decompositions of certain families of commuting partial differential operators.

Next, consider the following difference operator in the space of even functions depending on the variable $s$:

\begin{equation}
Lg(s) := \frac{(a - is)(b - is)}{(-2is)(-2is + 1)}(g(s + i) - g(s)) + \frac{(a + is)(b + is)}{(2is)(2is + 1)}(g(s - i) - g(s)),
\end{equation}

where $i^2 = -1$. A domain of this operator is a space of even functions holomorphic in the strip $|\text{Im } s| < 1 + \varepsilon$ with some condition of decreasing at infinity. It turns out that $L$ is essentially self-adjoint in the space of even functions $L^2_{\text{even}}(\mathbb{R}, |\Gamma(a + is)|\Gamma(b + is)|^2 ds)$ and the operator $I^{-1}_{a,b}$ sends it to the operator of multiplication by $x$.

So we have a bispectrality in the spirit of Grünbaum [17], [8]. Notice that simpler index integral transforms as the Kontorovich–Lebedev transform and the $1F_1$-Wimp transforms also are bispectral, see [38].

Cherednik showed [4] that inverse Heckman–Opdam transforms provide spectral decompositions of families of commuting difference operators, see also van Diejen, Emsiz [6].

1.3. Radial parts of Laplace operators. Recall one more classical topic. Consider the usual sphere $S^2_R$:

$$x^2 + y^2 + z^2 = 1,$$

the orthogonal group $\text{SO}(3)$ acts in $L^2(S^2_R)$ by rotations. Recall one of possible ways to decompose this unitary representation into irreducible components. Consider the Beltrami–Laplace operator $\Delta$ on the sphere and restrict it to the space of functions depending on the height $z$. We get a differential operator

$$L_z := (1 - z^2)\frac{\partial^2}{\partial z^2} - 2z\frac{\partial}{\partial z},$$

in $L^2[-1,1]$. Eigenfunctions of $L_z$ are the Legendre polynomials. Simple arguments show that the spectral decomposition of $\Delta$ is a priori equivalent to the spectral decomposition of $L_z$ (the reason of this equivalence is compactness of the group $\text{SO}(2)$ of rotations of $S^2_R$ about the vertical axis).

Now consider the complex manifold $S^2_C \subset \mathbb{C}^3$ defined by the same equation $x^2 + y^2 + z^2 = 1$. The complex orthogonal group $\text{SO}(3,\mathbb{C})$ (the Lorentz group) acts on the quadric $S^2_C$, the action admits an $\text{SO}(3,\mathbb{C})$-invariant measure, and again we come to a problem of decomposition of the unitary representation of $\text{SO}(3,\mathbb{C})$ in $L^2$ on $S^2_C$. Now we have two Beltrami-Laplace operators, a holomorphic operator $\Delta$ and an antiholomorphic operator $\overline{\Delta}$. They commute in the Nelson sense. Restricting them to functions depending on the coordinate $z \in \mathbb{C}$ we get two operators:

$$L_z := (1 - z^2)\frac{\partial^2}{\partial z^2} - 2z\frac{\partial}{\partial z}, \quad L_{\overline{z}} := (1 - \overline{z}^2)\frac{\partial^2}{\partial \overline{z}^2} - 2\overline{z}\frac{\partial}{\partial \overline{z}}.$$

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8This problem was solved by Naimark in [30] in a completely different way.

9This pair corresponds to $a = b = 1/2$ in our parameters.
However, now the stabilizer of the point \((x, y, z) = (0, 0, 1)\) is a *noncompact* subgroup \(SO(2, \mathbb{C})\), and this breaks the a priori argumentation. A joint spectral decomposition of \(\Delta, \overline{\Delta}\) can be reformulated as a certain problem\(^{10}\) for \(L_z, \overline{L}_z\), but this is not precisely a problem of a joint spectral decomposition of \(L_z, \overline{L}_z\).

Notice that a similar separation of variables can be done for \(L^2\) on an arbitrary rank one complex symmetric space \(G_C/H_C\) (and, more generally, for spaces of \(L^2\)-sections of line bundles on \(G_C/H_C\)). In all the cases we get pairs of hypergeometric operators of our type. We hope that our spectral decomposition allows to write the explicit Plancherel formula for such spaces and to give another proof of old Naimarks’s results \([30]–[32]\) on tensor products of representations of the Lorentz group. However, the present paper does not have such purposes.

### 1.4. Homographic transformations of the operators \(\mathfrak{D}, \overline{\mathfrak{D}}\).

Our next purpose is to present the explicit joint spectral decomposition of the pair \(\mathfrak{D}, \overline{\mathfrak{D}}\). We need some preparations.

Consider the following 8 transformations of functions on \(\hat{\mathbb{C}}\):

\[
\begin{align*}
  f(z) &\mapsto \gamma_j(z)f(z), & f(z) &\mapsto \gamma_j(z)f(1-z),
\end{align*}
\]

where

\[
\gamma_j(z) = 1, \quad |1-z|^{2(b-a)}, \quad |z|^{2(1-a-b)}, \quad |z|^{2(1-a-b)}|1-z|^{2(b-a)},
\]

cf. Erdélyi etc., \([9]\), Subsect. 2.6.1. It can be readily checked that these transformations send the operators \(\mathfrak{D}, \overline{\mathfrak{D}}\) to operators of the same type with other values of the parameters \((a, b)\), as

\[
(a, b) \mapsto (b, a), \quad (a, b) \mapsto (1-a, b), \quad \text{etc.}
\]

Thus we get all isometries of the square \(\Pi\). In particular, such transformations send spectral problems to equivalent spectral problems.

### 1.5. Notation. Generalized powers.

Denote by \(\mathbb{C}^\times\) the multiplicative group of \(\mathbb{C}\). We need a notation for characters of \(\mathbb{C}^\times\). Let \(z \in \mathbb{C}^\times\) and \(a, a' \in \mathbb{C}\) satisfy \(a-a' \in \mathbb{Z}\). We define a *generalized power* of \(z\) by

\[
z^a = z^{a|a'|} := z^{a \ln z + a' \ln z} = |z|^{2a|a'|-a},
\]

Denote by \(\Lambda_C\) the set of all pairs \(a|a'|\) such that \(a-a' \in \mathbb{Z}\). Denote by \(\Lambda \subset \Lambda_C\) the set of all pairs

\[
a|a' = \frac{1}{2}(k + is)\frac{1}{2}(-k + is), \quad \text{where } k \in \mathbb{Z}, \; s \in \mathbb{R}.
\]

Equivalently, \(a|a' \in \Lambda\) if \(a-a' \in \mathbb{Z}, \; a+a' \in i\mathbb{R}\). We also will use the notation

\[
[a] = [a|a'] := \frac{1}{2} \text{Re}(a + a').
\]

We have

\[
|z^a| = |z|^{2|a|},
\]

in particular, for \(a \in \Lambda\) we have \(|z^a| = 1\).

\(^{10}\) Such reductions for families of Laplace operators were widely used by Harish-Chandra (in his famous works on the Plancherel formula for real semisimple Lie groups) and by his successors. The problem for \(L_z, \overline{L}_z\) is more similar to decompositions of \(L^2\) on real rank one pseudo-Riemannian symmetric spaces, which was solved by one of the authors of the present paper \([27]–[29]\).
We fix the standard Lebesgue measure $d\lambda$ on the set $\Lambda$:
\[
\int_{\Lambda} \varphi(\lambda) \, d\lambda := \sum_k \int_{\mathbb{R}} \varphi\left(\frac{k+i\mu}{\sqrt{2}}\right) \, ds.
\]

1.6. Hypergeometric function of the complex field. Following Gelfand, Graev, and Retakh [33], we define the gamma function $\Gamma^C$, the beta function $B^C$, and the hypergeometric function $\text{}_2F_1^C$ of the complex field (see, also, Gelfand, Graev, Vilenkin, [34], Subsect. II.3.7, and Mimachi [26]). The gamma function $\Gamma^C$ is
\[
(1.16) \quad \Gamma^C(a) = \Gamma^C(a|a') := \frac{1}{\pi} \int_{\mathbb{C}} z^{a-1} e^{2i \text{Re}z} \, d\bar{z} := \frac{1}{\pi} \int_{\mathbb{C}} z^{a-1} e^{2i \text{Re}z} \, d\bar{z} = i^{a-a'} \frac{\Gamma(a)}{\Gamma(1-a')} = i^{a-a'} \frac{\Gamma(a')}{\Gamma(1-a')} = \frac{i^{a-a'} \Gamma(a')\sin \pi a'}{\pi}.
\]

The beta function $B^C$ is
\[
(1.17) \quad B^C(a,b) := \frac{1}{\pi} \int_{\mathbb{C}} (1-t)^{a-1} (1-b)\, d\bar{t} = \frac{\Gamma^C(a)\Gamma^C(b)}{\Gamma^C(a+b)} = \frac{\Gamma(a)\Gamma(b)\Gamma(1-a'-b')}{\Gamma(a+b)\Gamma(a')\Gamma(b')}.
\]

The hypergeometric function of the complex field is defined by
\[
(1.18) \quad \text{}_2F_1^C[a,b;c;z] := \text{}_2F_1^C[a,b;c;z] := \frac{1}{\pi B^C(b,c-b)} \int_{\mathbb{C}} t^{b-1}(1-t)^c(1-zt)^{-a} \, d\bar{t}.
\]

Recall that the Gauss hypergeometric functions are defined by
\[
\text{}_2F_1[a,b;c;z] := \frac{1}{B(b,c-b)} \int_0^1 t^{b-1}(1-t)^{c-b-1}(1-zt)^{-a} \, dt = \sum_{p=0}^{\infty} \frac{(a)_p(b)_p}{(c)_p} z^p,
\]
where $(c)_p := c(c+1)\ldots(c+p-1)$ is the Pochhammer symbol. The functions $\text{}_2F_1^C[a,b;c;z]$ admit expressions in the terms of $\text{}_2F_1$, see Theorem [59]

1.7. Spectral decomposition. For $(a,b) \in \Pi$ we define the kernel $K_{a,b}(z,\lambda)$ on $\mathbb{C} \times \Lambda$ by
\[
(1.19) \quad K_{a,b}(z,\lambda) = \frac{1}{\Gamma^C(a+b|a+b)} \text{}_2F_1^C\left[\frac{a+a'|a+b|a+b}{a+b|a+b};z\right].
\]

Theorem 1.3. Let $(a,b) \in \Pi_{\text{cont}}$. Then the operator
\[
J_{a,b}f(\lambda) := \int_{\mathbb{C}} K_{a,b}(z,\lambda) f(z) \, d\bar{z}
\]
is a unitary operator from $L^2(\mathbb{C},\mu_{a,b})$ to $L^2_{\text{even}}(\Lambda,\kappa_{a,b})$ of even functions on $\Lambda$ with respect to the Plancherel measure
\[
(1.20) \quad dK_{a,b}(\lambda) = \kappa_{a,b}(\lambda) \, d\lambda = \frac{1}{4\pi^2} \left|\lambda \Gamma^C(a-\lambda|a+b) \Gamma^C(b+\lambda|b-a)\right|^2 \, d\lambda.
\]

\[^{11}\text{This integral has a multi-dimensional counterpart of the Selberg type, see [7].}\]
Next, we modify the definition of the measure for \((a, b) \in \Pi \setminus \Pi_{\text{cont}}\). Due to the homographic transformations\(^\text{12}\) it is sufficient to consider the case \(a < 0\). We define the Plancherel measure \(dK_{a,b}(\lambda)\) on \(\Lambda\) that is the sum of \(\nu_{a,b} \, d\lambda\) and two \(\delta\)-measures located at the points \(\pm a\) \(\in \Lambda\).

\[
\Gamma^C(a + b|a + b) \Gamma^C(b - a|b - a) \Gamma^C(2a|2a) \cdot (\delta_{a[a + \delta - a]} - \delta_{a[a]}).
\]

Define a constant function \(v(z)\) on \(\hat{\mathbb{C}}\) by

\[
v(z) = \Gamma^C(a + b|a + b)^{-1}.
\]

For \(f \in D(\hat{\mathbb{C}})\) we define an even function \(J_{a,b}(\lambda)\) on the support of \(dK_{a,b}(\lambda)\) given by the same formula\(^\text{12}\) on \(\Lambda\), its value at \((\pm a, \pm a)\) is

\[
J_{a,b}f(\pm a|\pm a) := (f, v)_{L^2(\mathbb{C}, \mu_{a,b})}.
\]

**Theorem 1.4.** If \((a, b) \in \Pi\) and \(a < 0\), then the operator \(J_{a,b}\) is unitary as an operator \(L^2(\mathbb{C}, \mu_{a,b})\) to \(L^2_{\text{even}}(\Lambda, dK_{a,b})\).

Our operator really determines the spectral decomposition:

**Theorem 1.5.** For each \((a, b) \in \Pi\) for any \(f \in D_{\text{even}}(\hat{\mathbb{C}})\) we have

\[
J_{a,b} \mathcal{D} f(\lambda) = \lambda^2 J_{a,b} f(\lambda), \quad J_{a,b} \mathcal{F} f(\lambda) = \mathcal{F}^2 J_{a,b} f(\lambda).
\]

This means that

\[
\mathcal{D}\mathcal{K}(z, \lambda) = \lambda^2 \mathcal{K}(z, \lambda), \quad \mathcal{F}\mathcal{K}(z, \lambda) = \mathcal{K}^2 \mathcal{K}(z, \lambda).
\]

Next, we consider the space \(D_{\text{even}}(\hat{\mathbb{A}})\), which consists of even smooth compactly supported functions on \(\Lambda\) that are zero on a neighborhood of the point \(0|0\). The following statement explains the appearance of the space \(\mathcal{R}_{a,b}\) and also is one of the arguments for the proofs of our main statements.

**Theorem 1.6.** If \(F \in D_{\text{even}}(\hat{\mathbb{A}})\), then \(J_{a,b}^* F \in \mathcal{R}_{a,b}\).

The images of \(\delta\)-functions also are contained in \(\mathcal{R}_{a,b}\).

### 1.8. The transformation \(J_{a,b}\) in the complex domain.

Let us extend our kernel \(\mathcal{K}\) to the complex domain. For

\[
\{\lambda|\lambda'\} = \left\{\frac{k + \sigma}{2}, \frac{-k + \sigma}{2}\right\} \in \Lambda
\]

we set

\[
\mathcal{K}_{a,b}(z; \lambda) = \mathcal{K}_{a,b}(z; k, \sigma) := \frac{1}{\Gamma^C(a + b|a + b)} 2 F_1 \left[ a + \frac{k + \sigma}{2}, a + \frac{-k + \sigma}{2}; \frac{a + k - \sigma}{a + b}; \frac{a + \frac{k - \sigma}{2}}{a + b}; z \right],
\]

where \(k\) ranges in \(\mathbb{Z}\), \(\sigma\) ranges in \(\mathbb{C}\). The previous expression\(^\text{12}\) corresponds to a pure imaginary \(\sigma\).

For \(f \in D(\hat{\mathbb{C}})\) we define a meromorphic function on \(\Lambda\) by

\[
J_{a,b}(k, \sigma) := \int_{\hat{\mathbb{C}}} f(z) \mathcal{K}(z; k, \sigma) \, d\mu_{a,b}(z).
\]

**Theorem 1.7.** For \(f \in D(\hat{\mathbb{C}})\) the function \(J_{a,b} f\) is contained in the space \(\mathcal{W}_{a,b}\) defined as follows.

\(^{12}\) Changing of kernels \(\mathcal{K}_{a,b}\) by the homographic transformation can be observed from Proposition \(\text{6.6}\).
We define a space \( W_{a,b} \) as the space of all meromorphic functions \( F(k, \sigma) \) on \( \Lambda_C \) satisfying the conditions a)–d):

a) \( F \) is even, i.e., \( F(-k, -\sigma) = F(k, \sigma) \).

b) Possible poles of \( F(k, \sigma) \) are located at the points

\[
\sigma = \pm (-2a + |k| + 2j), \quad \pm (-2b + |k| + 2j), \quad \text{where } j = 1, 2, 3, \ldots
\]

A maximal possible order of a pole at a point \((l, c)\) is a multiplicity of \((l, c)\) in the collection \(^{14}[1,13]\).

c) For each \( A > 0 \) for each \( N > 0 \) in the union of strips \(|\text{Re}\,\sigma| < A\) we have an estimate

\[
F(k, \sigma) = O(k^2 + (\text{Im}\,\sigma)^2)^{-N} \quad \text{as } k^2 + (\text{Im}\,\sigma)^2 \to \infty.
\]

d) For each \( p, q \in \mathbb{Z} \)

\[
F(p, q) = F(q, p).
\]

Next, we extend the spectral density \( \kappa_{a,b} \) to the complex domain.

\[
\kappa_{a,b}(\lambda|\lambda') = \kappa_{a,b}(k, \sigma) :=
\]

\[
\frac{1}{4\pi^2} (k + \sigma)(k - \sigma) \Gamma^C(a + \frac{k + \sigma}{2}) \Gamma^C(a + \frac{-k - \sigma}{2}) \Gamma^C(b + \frac{k + \sigma}{2}) \Gamma^C(b + \frac{-k - \sigma}{2}).
\]

In the case \( a < 0 \) discussed above, \( \kappa_{a,b} \) has a pole at \( k = 0, \sigma = a \) and the inner product in \( L^2_{\text{even}}(\Lambda, dK_{a,b}) \) can be written as

\[
\langle F, G \rangle = \frac{1}{i} \sum_k \int_{-i\infty}^{i\infty} F(k, \sigma) \overline{G(k, -\sigma)} \kappa_{a,b}(k, \sigma) \, d\sigma + 2 \text{res}_{s=a} \left( F(k, \sigma) \overline{G(k, -\sigma)} \kappa_{a,b}(0, \sigma) \right) \).
\]

If \( a > 1 \), then the spectral density has a zero at \( k = 0, \sigma = a - 1 \) but both functions \( F(k, \sigma), \overline{G(k, -\sigma)} \) admit simple poles at this point, and we have a similar formula.

**1.9. Difference spectral problem.** It turns out that our problem is bispectral, and the bispectrality is a crucial argument of our proof. We define analogs of the difference operator \(^{13}[1,13]\). Consider meromorphic functions \( \Phi \) depending on

\[
\lambda|\lambda' = \frac{1}{2}(k + is)|\frac{1}{2}(-k + is) \in \Lambda_C
\]

and the operators in the space of meromorphic functions defined by

\[
T\Phi(k, s) = \Phi(k + 1, s - i), \quad \overline{T}\Phi(k, s) = \Phi(k + 1, s + i),
\]

or, equivalently,

\[
T\Phi(\lambda|\lambda') = \Phi(\lambda + 1|\lambda'), \quad \overline{T}\Phi(\lambda|\lambda') = \Phi(\lambda|\lambda' + 1).
\]

\(^{13}\)We say that a function \( F(k, \sigma) \) is meromorphic if it is meromorphic as a function in \( \sigma \) for any fixed \( k \).

\(^{14}\)For \((a, b) \in \Pi\) orders of poles \( \leq 2 \). Poles of order 2 arise only if \( a = b, a = 1, b = 1 \).
We define the following difference operators

\begin{align}
(1.28) \quad \mathcal{L} := & \frac{(a + \lambda)(b + \lambda)}{2\lambda(1 + 2\lambda)} (T - 1) + \frac{(a - \lambda)(b - \lambda)}{-2\lambda(1 - 2\lambda)} (T^{-1} - 1); \\
(1.29) \quad \mathcal{T} := & \frac{(a + \lambda')(b + \lambda')}{2\lambda'(1 + 2\lambda')} (\widetilde{T} - 1) + \frac{(a - \lambda')(b - \lambda')}{-2\lambda'(1 - 2\lambda')} (\widetilde{T}^{-1} - 1).
\end{align}

Formally,

\[ \mathcal{L}\mathcal{T} = \mathcal{T}\mathcal{L}. \]

**Theorem 1.8.** a) The operators \( \frac{1}{2}(\mathcal{L} + \mathcal{T}) \), \( \frac{1}{2}(\mathcal{L} - \mathcal{T}) \) defined on the space \( \mathcal{W}_{a,b} \) are essentially self-adjoint and commute in the Nelson sense.

b) For \( \Phi \in J_{a,b}\mathcal{D}(\mathcal{C}) \) we have

\begin{equation}
(1.30) \quad J_{a,b}^{-1} \mathcal{L} \Phi(z) = z J_{a,b}^{-1} \Phi(z), \quad J_{a,b}^{-1} \mathcal{T} \Phi(z) = z J_{a,b}^{-1} \Phi(z).
\end{equation}

Thus the operator \( J_{a,b}^{-1} \) determines a joint spectral decomposition of \( \frac{1}{2}(\mathcal{L} + \mathcal{T}) \) and \( \frac{1}{2}(\mathcal{L} - \mathcal{T}) \).

**1.10. The structure of the proofs.** We derive asymptotics of the kernel \( \mathcal{K}(z, \lambda) \) as \( z \to 0, 1, \infty \) for fixed \( \lambda \) (Theorem 3.9) and as \( |\lambda| \to \infty \) for fixed \( z \) (Theorem 1.11). Next, we prove inclusions

\[ J_{a,b}^* \mathcal{D}_{\text{even}}(\mathcal{A}) \subset \mathcal{R}_{a,b}, \quad J_{a,b} \mathcal{D}(\mathcal{C}) \subset \mathcal{W}_{a,b} \]

(Proposition 5.2 and Corollary 8.2) and symmetries

\begin{align}
(1.31) \quad &\langle \mathcal{D} f, g \rangle_{L^2(\mathbb{C}, \mu_{a,b})} = \langle f, \mathcal{T} g \rangle_{L^2(\mathbb{C}, \mu_{a,b})}, \quad \text{where } f, g \in \mathcal{R}_{a,b}; \\
(1.32) \quad &\langle \mathcal{E} F, G \rangle_{L^2(\Lambda, dK_{a,b})} = \langle F, \mathcal{T} G \rangle_{L^2(\Lambda, dK_{a,b})}, \quad \text{where } F, G \in \mathcal{W}_{a,b},
\end{align}

see Proposition 5.5 and Theorem 8.4. This implies a generalized orthogonality, i.e.,

\[ \langle J_{a,b}^* F, J_{a,b}^* G \rangle_{L^2(\mathbb{C}, \mu_{a,b})} = 0 \]

if supports of \( F, G \in \mathcal{D}_{\text{even}}(\mathcal{A}) \) are disjoint, and a similar statement for \( J_{a,b} F, G \), see Lemmas 6.3, 6.4. Next, we show that for any \( F, G \in \mathcal{D}_{\text{even}}(\mathcal{A}) \) the inner products of their preimages can be written as:

\[ \langle J_{a,b}^* F, J_{a,b}^* G \rangle_{L^2(\mathbb{C}, \mu_{a,b})} = \langle F, G \rangle_{L^2(\Lambda, dK_{a,b})} + \int \int H(\lambda_1, \lambda_2) F(\lambda_1) G(\lambda_2) d\lambda_1 d\lambda_2,
\]

where \( H \) is a locally integrable function, see Lemma 6.4. We also prove a similar statement for \( J_{a,b} F, G \), see Lemma 6.4. Then generalized orthogonality implies \( H(\cdot, \cdot) = 0 \). Thus we get

\begin{equation}
(1.33) \quad J_{a,b}^* J_{a,b} = 1, \quad J_{a,b} J_{a,b} = 1,
\end{equation}

and this is our main statement.

Some steps of this double way are straightforward, some points require long calculations and estimates, and we meet some points of good luck (the proofs of Theorem 8.4 and Lemma 9.3). We also need a lot of information about functions \( \left( \begin{array}{c} 2 F_1^C \\ 2 F_1^C \end{array} \right) \) (in particular, to cover the cases \( a + b \in \mathbb{Z} \) and \( a - b \in \mathbb{Z} \) we need a tedious examination of possible degenerations of functions \( \left( \begin{array}{c} 2 F_1^C \\ 2 F_1^C \end{array} \right) \).

The bispectrality allows to avoid a direct proof of completeness of the system of generalized eigenfunctions of \( \mathcal{D}, \mathcal{D} \).

To prove the necessary conditions of self-adjointness in Theorem 1.11 we analyze common generalized eigenfunctions of \( \mathcal{D}, \mathcal{D} \) for \( (a, b) \notin \Pi \) and after a natural
selection we reduce a set of possible candidates to a finite family. This is done in Section 4.

This text is focused to a proof of unitarity of $J_{a,b}$ in Section 4. An introduction to functions $2F_1^C$ in Section 3 can be a point of an independent interest. Also, we get two relatively pleasant statements about asymptotic behavior of integrals

$$M(\varepsilon) = \int_{\mathbb{C}} t^{\alpha-1}|\alpha'|^{-1}(\varepsilon - t)^{\beta-1}|\beta'|^{-1}\psi(t) \, d\overline{\overline{t}} \quad \text{as } \varepsilon \to 0$$

and

$$I(\lambda) = \int_{\mathbb{C}} |f(t)|^2 e^{i \Re(\lambda \varphi(t))} \, d\overline{\overline{t}} \quad \text{as } |\lambda| \to \infty,$$

where $f, \varphi$ are holomorphic and $\lambda \in \mathbb{C}$ (Theorems 2.3 and 7.2).

1.11. Final remarks. The index hypergeometric transform (1.11) can be applied as a heavy tool of theory of special functions, see [22], [35], [37]. In [39] we use our operators $J_{a,b}$ to obtain a beta integral over $\Lambda$, which is a counterpart of the Dougall $5H_5$-summation formula and of the de Branges–Wilson integral.

Also, we notice that functions, which can be regarded as higher hypergeometric functions $4F_3^C$ of the complex field, arise in a natural way in the work of Ismagilov [20] as analogs of the Racah coefficients for unitary representations of the Lorentz group $\text{SL}(2, \mathbb{C})$ (see, also a continuation in [5]).

It seems that our problem can be a representative of some family of spectral problems, but now it is too early to claim something certainly.

2. Preliminaries. Gamma function, the Mellin transform, weak singularities

This section is a union of 3 disjoint topics:

— some properties of the function $\Gamma^C$, which are intensively used below:

— some properties of the Mellin transform on $\mathbb{C}$, they are used in a proof of Proposition 3.1 and in Sections 7–9:

— a lemma from asymptotic analysis, which is used only in a proof of Theorem 7.9 (the last statement can be independently established by a straightforward tiresome way):

2.1. Some properties of the gamma function. The usual functional equations for the $\Gamma$-function can be easily rewritten for $\Gamma^C$ (recall that $a - a' \in \mathbb{Z}$!):

$$\Gamma^C(a|a') = \Gamma^C(a'|a);$$  

$$\Gamma^C(a + 1|a') = i a \Gamma^C(a|a');$$  

$$\Gamma^C(a|a') \Gamma^C(1 - a|1 - a') = (-1)^{a-a'};$$  

$$\Gamma^C(a|a') = (-1)^{a-a'} \Gamma^C(\overline{a}|\overline{a'}).$$

Also,

$$\prod_{p=0}^{m-1} \Gamma^C \left( a + \frac{p-1}{m}, |a'| + \frac{p-1}{m} \right) = m^{1-m(a+a')} \Gamma^C(m|ma').$$

The identity (2.4) implies

$$B^C[a|a', b|b'] = B^C[\overline{a}|\overline{a'}, \overline{b}|\overline{b'}].$$
Let $k_1, k_2 \in \mathbb{Z}$. Then

$$\Gamma_C(k_1|k_2) = \begin{cases} \infty, & \text{if } k_1, k_2 \in \mathbb{Z}_-, \\ 0, & \text{if } k_1, k_2 \in \mathbb{N}, \\ i^{k_1-k_2} \frac{(k_1-1)!}{(k_2-1)!}, & \text{if } k_1 \in \mathbb{N}, k_2 \in \mathbb{Z}_-, \\ i^{k_2-k_1} \frac{(k_2-1)!}{(k_1-1)!}, & \text{if } k_2 \in \mathbb{N}, k_1 \in \mathbb{Z}_-, \end{cases}$$

(2.6)

where $\mathbb{Z}_-$ denotes the set of integers $\leq 0$.

The following lemma gives us the asymptotics of the Plancherel density (1.20).

**Lemma 2.1.** We have the following asymptotics in $\lambda \in \Lambda$:

$$\Gamma_C(a - \lambda a + \bar{\lambda}) \Gamma_C(b + \lambda |b - \bar{\lambda}) \sim \lambda^{a+b-1} |a+b-1| \quad \text{as } |\lambda| \to \infty.$$  

The asymptotics is uniform in $a, b$ if they range in a bounded domain.

**Proof.** Denote $\Re \lambda = k/2$. Let $|\arg \lambda| < \pi - \varepsilon$. Then we write our expression as

$$i^k \Gamma(a + \bar{\lambda}) / \Gamma(1-a+\lambda)$$

and apply the standard asymptotic formula $\Gamma(z+\alpha)/\Gamma(z+\beta) \sim z^{\alpha-\beta}$ in the sector $|\arg z| < \pi - \varepsilon$, see Erdélyi etc., [9], formula (1.18.4). If $|\arg(-\lambda)| < \pi - \varepsilon$, we write

$$i^{-k} \Gamma(1-a-\lambda) / \Gamma(1-a+\lambda),$$

and come to the same asymptotics. \qed

**2.2. The Mellin transform.** Denote by $\mathbb{C}^\times := \mathbb{C} \setminus \{0\}$ the multiplicative group of $\mathbb{C}$. The Mellin transform (see, e.g., [13]) on $\mathbb{C}^\times$ is defined by

$$g(\mu) = \mathcal{M}f(\mu) = \frac{1}{2\pi} \int_{\mathbb{C}} f(z) z^{-\mu-1} d\bar{z},$$

(2.8)

where $\mu = \{\mu | \mu'\} = \{\frac{k+i\varepsilon}{2}, \frac{k+i\varepsilon}{2}\} \in \Lambda_\mathbb{C}$ (here we allow complex $s$). This operator is the Fourier transform on the group $\mathbb{C}^\times \simeq (\mathbb{R}/2\pi \mathbb{Z}) \times \mathbb{R}$, so it is reduced to the Fourier transforms on $(\mathbb{R}/2\pi \mathbb{Z})$ and on $\mathbb{R}$. Indeed, changing variables

$$z = e^{\rho} e^{i\varphi},$$

we come to

$$g(k, s) = \frac{1}{2\pi} \int_0^{2\pi} \int_{-\infty}^{\infty} f(e^{\rho} e^{i\varphi}) e^{ik \varphi + is \rho} d\rho d\varphi.$$  

The inversion formula is given by

$$f(z) = \mathcal{M}^{-1} g(\mu) = \frac{1}{2\pi} \int_{\Lambda} g(\mu | -\mu) z^{-\mu} d\mu.$$  

Equivalently, $\mathcal{M}$ is a unitary operator $L^2(\mathbb{C}^\times, |z|^{-2}) \to L^2(\Lambda)$.

**2.3. The Mellin transform of even functions.** We say that a function $f$ on $\mathbb{C}^\times$ is $\times$-even if $f(z^{-1}) = f(z)$. Denote by $L^2_\times(\mathbb{C}^\times, |z|^{-2})$ the subspace of $L^2(\mathbb{C}^\times, |z|^{-2})$ consisting of $\times$-even functions. Obviously, the Mellin transform sends $\times$-even functions in $z$ to even functions in $\mu$. Also, for a $\times$-even function $f$ we have

$$\mathcal{M}f(\mu) = \frac{1}{2} \int_{\mathbb{C}} f(z) (z^{\mu-1} + z^{-\mu-1}) d\bar{z},$$

(2.9)

where $f$ is $\times$-even.

**2.4. The Mellin transform of smooth compactly supported functions.**
Theorem 2.2. a) Let $f$ be a compactly supported smooth function on $\mathbb{C}$. Then $\mathcal{M}f(\mu|\mu')$ extends to a meromorphic function in the variable $\mu$ with possible poles located at the points $\mu|\mu' \in \mathbb{Z}_- \times \mathbb{Z}_-$. Moreover, for any $p, p' \in \mathbb{Z}_+$ for $\text{Re}(\mu + \mu') > -p - p'$ we have

$$I(\mu|\mu') = \frac{(-1)^{p+p'}}{2\pi(-\mu)_p (-\mu')_{p'}} \int_{\mathbb{C}} z^{\mu-p-1|\mu'-p'-1} \frac{\partial^{p+p'}}{\partial z^p \partial z'^{p'}} f(z) \, dz.$$  

The residues at the poles are

$$\text{res}_{\mu|\mu'=-p|-p'} I(\mu|\mu') = \frac{1}{(p-1)! (p'-1)!} \frac{\partial^{p+p'}}{\partial z^p \partial z'^{p'}} f(0,0).$$

b) For each $N$ for each $A$ for all pairs $(k,s)$ satisfying $|\text{Im} s| < A$ we have

$$\mathcal{M}f\left(\frac{k+is}{2} | -\frac{k+is}{2}\right) = O(k^2 + |s|^2)^{-N} \quad \text{as } |k^2| + |s^2| \to \infty.$$ 

For a proof of statement a), see Gelfand, Shilov [15], Sect. B.1.3, or equivalently Russian edition of Gelfand, Graev, Vilenkin [14], Addendum 1.3 (the term 'Mellin transform' in that place is absent, but the statement is proved). Formula (2.10) is obtained from (2.8) by integration by parts. The statements about location of poles and about residues require more careful arguments.

**Proof of statement b).** We pass to polar coordinates, $z = e^{i\theta}$ and get

$$\mathcal{M}f\left(\frac{k+is}{2} | -\frac{k+is}{2}\right) = \frac{1}{2\pi} \int_0^\infty \int_0^{2\pi} H(\theta, r) r^{-1+ik} e^{i\theta k} \, d\theta \, dr,$$

where $H(\theta, r) := \Phi(r e^{i\theta})$ is a smooth function $2\pi$-periodic in $\theta$, the $H(\theta, 0)$ does not depend on $\theta$, also $H(\theta + \pi, -r) = H(\theta, r)$. Integrating by parts in $r$, we get

$$\mathcal{M}f\left(\frac{k+is}{2} | -\frac{k+is}{2}\right) = \frac{(-1)^l}{2\pi (ik)^m(s)l} \int_0^{2\pi} \int_0^\infty \frac{\partial^l}{\partial r^l} H(\theta, r) r^{-1+is+l} \, dr \, e^{i\theta k} \, d\theta.$$ 

For $l > A$ the integral absolutely converges. Integrating by parts in $\theta$, we get

$$\frac{(-1)^l}{2\pi (ik)^m(s)l} \int_0^{2\pi} \int_0^\infty \frac{\partial^{l+m}}{\partial r^l \partial \theta^m} H(\theta, r) r^{-1+is+l} \, dr \, e^{i\theta k} \, d\theta,$$

and

$$|\mathcal{M}f\left(\frac{k+is}{2} | -\frac{k+is}{2}\right)| \leq \frac{\text{const}}{|(2\pi ik)^m(s)l|}.$$ 

If $|s| > |k|$ we take $m = 0$ and large $l$, if $|k| > |s|$, we take $l > |\text{Im} s|$ and large $m$. $\square$

2.5. **Weak singularities.** Here we imitate one standard trick of asymptotic analysis, see, e.g., [10], Sect. 1.4. Fix $R$ and a smooth function $\psi(t)$ on $\mathbb{C}$. Consider the integrals of the following type

$$M(\varepsilon) = M_{\alpha,\beta}(\varepsilon) := \int_{|t| \leq R} t^{\alpha-1|\alpha'|-1}(\varepsilon - t)^{\beta-1|\beta'|-1} \psi(t) \, d\bar{t}.$$ 

Clearly, $M_{\alpha,\beta}(\varepsilon)$ is holomorphic in $\alpha, \beta$ in the domain of convergence and admits a meromorphic continuation to $(\alpha, \beta) \in \Lambda^2$.  

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15 For instance, see the proof of Proposition [14] below.
Theorem 2.3. Let \( \alpha, \beta \) satisfy the condition
\begin{equation}
\alpha, \beta, \alpha + \beta - 1 \notin \mathbb{Z}_- \times \mathbb{Z}_-
\end{equation}
Then \( M(\varepsilon) \) (defined in the sense of analytic continuation) admits the following asymptotic expansion at 0:
\begin{equation}
M(\varepsilon) \sim \sum_{i, j' \geq 0} B^C(\alpha + i |\alpha'| + i', \beta |\beta'|) \cdot \frac{1}{i! j'!} \frac{\partial^{i + j'} \psi(0, 0)}{\partial t^i \partial \bar{t}^{j'}} \cdot \varepsilon^{\alpha + \beta + i - 1 |\alpha'| + \beta' + i' - 1 +}
+ \sum_{j, j' \geq 0} r_{j|j'} \varepsilon^{|j||j'|}.
\end{equation}
The coefficients of the expansion are meromorphic in \( \alpha |\alpha'|, \beta |\beta'| \).
\( R^{[\varepsilon]} \) if \( |\alpha| > 0, |\beta| > 0, |\alpha| + |\beta| > 1 \), then
\begin{equation}
r_{00} = \int_{|t| < R} t^{\alpha + \beta - 2 |\alpha'| + \beta' - 2} \psi(t) \, dt.
\end{equation}

First, we prove the following lemma

Lemma 2.4. Let \( \alpha, \alpha', \beta, \beta' \) satisfy the condition (2.12). Then the following integral (defined in the sense of analytic continuation)
\begin{equation}
R(\varepsilon) = \int_{|t| < R} t^{\alpha - 1 |\alpha'| - 1} \varepsilon^{\beta - 1 |\beta'|- 1} \, dt
\end{equation}
admits an asymptotic expansion of the form
\begin{equation}
R(\varepsilon) = B^C(\alpha |\alpha'|, \beta |\beta'|) \cdot \varepsilon^{\alpha + \beta - 1 |\alpha'| + \beta' - 1} + \sum_{j \geq 0, j' \geq 0} p_{j|j'} \varepsilon^{|j||j'|}.
\end{equation}
Moreover, the series
\begin{equation}
\sum_{j \geq 0, j' \geq 0} p_{j|j'} \varepsilon^{|j||j'|}
\end{equation}
converges in the circle \( |\varepsilon| < 1/R \), and the coefficients \( p_{j|j'}(a, b) \) are holomorphic in the domain (2.12).

PROOF. Set
\[
Q_{\alpha, \beta}(\varepsilon_1, \varepsilon_2) := (-1)^{j' - \beta} \int_{|t| > R} t^{\alpha + \beta - 2 |\alpha'| + \beta' - 2} \left(1 - \frac{\varepsilon_1}{t}\right)^{\beta - 1} \left(1 - \frac{\varepsilon_2}{t}\right)^{\beta' - 1} \, dt.
\]
This function is meromorphic in \( \alpha, \beta \), and in \( \varepsilon_1, \varepsilon_2 \) in the bidisk \( |\varepsilon_1| < 1/R, |\varepsilon_2| < 1/R \). Let
\begin{equation}
|\alpha| > 0, |\beta| > 0, |\alpha| + |\beta| < 1.
\end{equation}
Under these conditions the integral \( R(\varepsilon) \) converges, and
\[
R(\varepsilon) = \int_{|t| > R} B^C(\alpha |\alpha'|, \beta |\beta'|) \cdot \varepsilon^{\alpha + \beta - 1 |\alpha'| + \beta' - 1} - Q_{\alpha, \beta}(\varepsilon, \bar{\varepsilon}).
\]
\(^{16}\)Recall notation (1.15).
Expanding the integrand in \( Q_{\alpha, \beta} \) in a series in \( \varepsilon_1, \varepsilon_2 \) and integrating termwise we come to

\[
(2.19) \quad Q_{\alpha, \beta}(\varepsilon_1, \varepsilon_2) = 
\sum_{j \geq 0, j' \geq 0: \alpha + \beta - j = \alpha' + \beta' - j'} \frac{(-\beta + 1)_j (-\beta' + 1)_{j'}}{(j + j' - \alpha - \alpha' - \beta - \beta') j! j'} \varepsilon_1^j \varepsilon_2^{j'}.
\]

Now we can omit restrictions (2.18). Indeed, under conditions (2.12) the series (2.19) converges in the bidisk \(|\varepsilon_1| < 1, |\varepsilon_2| < 1\) and therefore its sum coincides with the meromorphic continuation.

**Proof of Theorem 2.3.** We expand the function \( \psi \) as a sum

\[
\psi(t, \bar{t}) = \sum_{j + j' \leq N} \frac{t^j \bar{t}^{j'}}{j! j'} \psi(0) t^{j'} H_N(t),
\]

where \( H_N(t) \) is a smooth function and

\[
H_N(t) = O(|t|^{N+1}) \quad \text{as } t \to 0.
\]

Substituting this to the initial integral we get a sum of integrals of the form (2.15), we apply Lemma 2.4 to each summand. Also we get a summand

\[
I(\varepsilon) = \int_{|t| \leq R} t^{\alpha-1} |\varepsilon + t|^{\beta-1} |\varepsilon - t|^{\beta'-1} H_N(t) d\bar{t}.
\]

We wish to show that \( T(\varepsilon) \) has partial derivatives at 0 up to order \( N - k \), where \( k \) is constant depending only on \( \alpha \) and \( \beta \). Consider a partition of unity, \( 1 = \varphi_1 + \varphi_2 \) such that \( \varphi_2 \) is zero at some smaller circle \(|t| < R'\). According to this, we split \( I = I_1 + I_2 \). Obviously, \( I_2 \) has an expansion of the form

\[
I_2 \sim \sum_{j, j' \geq 0} c_{j j'} |\varepsilon|^{j} j'!
\]

with coefficients meromorphic in \( \alpha, \beta \). Next, we integrate \( I_1 \) by parts several times,

\[
I_1(\varepsilon) = \frac{1}{(\beta)_m (\bar{\beta})_m} \int_{|t| \leq R} (\varepsilon + t)^{\beta - 1 + m} |\beta' + m - 1| H_N(t) \varphi_1(t) d\bar{t}.
\]

Choosing \( m \) we can make \( \beta + m - 1, \beta' + m - 1 \) as large, as we want, say > \( q \). Next, we choose a large \( N \), such that \( \frac{\partial^{2m}}{\partial \varepsilon \partial \bar{t}} (\ldots) \) is continuous at 0. Now we can differentiate \( I_1(\varepsilon) \) with respect to \( \varepsilon, \bar{\varepsilon} \) \( q \) times at 0 and consider its Taylor expansion. This finishes a derivation of the asymptotic expansion for \( R(\varepsilon) \).

If the integral \( R(0) \) converges, we substitute \( \varepsilon = 0 \) to the expansion and get the expression for \( r_{00} \).

**3. The hypergeometric function of the complex field**

Here we discuss basic properties of the functions \( _2F_1^C[\cdot] \).

**3.1. Domain of convergence and analytic continuation.** The hypergeometric function \( _2F_1^C[a, b; c; z] \) of the complex field is defined by the Euler type integral [1, 18]:

\[
(3.1) \quad _2F_1^C[a, b; c; z] = \frac{1}{\pi B^C(b, c - b)} \int_C t^{b-1} (1 - t)^{c-b-1} (1 - zt)^{-a} d\bar{t}.
\]
For \( z \neq 0, 1 \), the integral absolutely converges (see notation (1.15) if \( a, b, c \) is contained in the following tube \( \hat{\Xi} \),
\[
\hat{\Xi} : \quad \begin{cases} |b| > 0, & |c - b| > 0, \quad |a| < 1, \\ |c - a| < 2. \end{cases}
\]
In other words, the integral absolutely converges if and only if the point \(([a], [b], [c])\) is contained in the simplex \( \Xi \) in \( \mathbb{R}^3 \) with vertices \((1,0,0), (-1,0,0), (1,0,1), (1,2,2)\).

We have \( \Lambda^C \simeq \mathbb{C} \times \mathbb{Z} \), therefore triples \((a, b, c)\) depend on 3 integers and 3 complex parameters. Clearly, each component of the set \( \mathbb{Z}^3 \times \mathbb{C}^3 \) has an open intersection with the domain of convergence\(^{17}\).

**Proposition 3.1.** For \( z \in \hat{\mathcal{C}} \), the expression \( 2F^C_1[a, b; c; z] \) as a function of \( a, b, c \) admits a meromorphic extension to arbitrary values of \( a, b, c \) with poles at a countable union of surfaces
\[
(3.4) \quad a \in \mathbb{N} \times \mathbb{N}, \quad b \in \mathbb{N} \times \mathbb{N}, \quad c - a \in \mathbb{N} \times \mathbb{N}, \quad c - b \in \mathbb{N} \times \mathbb{N},
\]
and vanishes for all \( z \in \hat{\mathcal{C}} \) at
\[
(3.5) \quad c \in \mathbb{Z}_- \times \mathbb{Z}_-.
\]

**Proof.** Consider a partition of unity \( 1 = \varphi_0(t) + \varphi_1(t) + \varphi_{1/2}(t) + \varphi_\infty(t) + \varphi_\varnothing(z) \), where all summands are smooth and nonnegative, \( \varphi_0, \varphi_1, \varphi_{1/2}, \varphi_\infty \) are zero outside small neighborhoods of \( 0, 1, 1/2, \infty \), respectively, and \( \varphi_\varnothing = 0 \) in small neighborhoods of these points. Denote \( P(t, \tilde{t}) \) the integrand in the integral representation of \( 2F^C_1[a, b; c; z] \). Then
\[
\pi B^C(b, c - b) \cdot 2F^C_1[a, b; c; z] =
\]
\[
= \int \varphi_0 P \, d\tilde{t} + \int \varphi_1 P \, d\tilde{t} + \int \varphi_{1/2} P \, d\tilde{t} + \int \varphi_\infty P \, d\tilde{t} + \int \varphi_\varnothing P \, d\tilde{t}.
\]
The last summand is an entire function in \( a, b, c \). By Theorem 2.2 other summands are meromorphic and can have poles at
\[
b \in \mathbb{Z}_- \times \mathbb{Z}_-, \quad c - b \in \mathbb{Z}_- \times \mathbb{Z}_-, \quad a \in \mathbb{N} \times \mathbb{N}, \quad c - a \in \mathbb{N} \times \mathbb{N}.
\]
However, \( B^C \)-factor in the front of the integral \((3.1)\) kills the first and the second families of poles and produces new poles and also zeros. This gives us \((3.3)-(3.6)\), in particular the factor \( \Gamma^C(c) \) produces poles \((3.5)-(3.6)\).

All these possible poles really are poles, the simplest way to observe this is to look at formulas \((3.26)-(3.34)\) derived below. Formulas \((3.26)-(3.28)\) show that \((3.4)\) are poles. To check a presence of poles \((3.5)-(3.6)\) we apply \((3.8)-(3.11). \boxdot

**3.2. Kummer symmetries.** This section contains a collection of elementary formulas, they partially depend on Theorem 3.9 proved below. However, our proof of this theorem is based on differential equations and asymptotic analysis and is independent of our formulas.

\(^{17}\)The map \((a, b, c) \rightarrow ([a], [b], [c])\) from \( \Lambda^3 \rightarrow \mathbb{R}^3 \) is surjective on all components.
First we notice two trivial identities

\( (3.7) \quad \begin{aligned} 2F_1^c \left[ \frac{a|a', b|b'}{c' | c} ; z \right] &= 2F_1^c \left[ \frac{a' | a, b'}{c' | c} ; z \right] ; \end{aligned} \)

\( (3.8) \quad \begin{aligned} 2F_1^c \left[ \frac{a' | a, b'}{c' | c} ; z \right] &= 2F_1^c \left[ \frac{a' | a, b'}{c' | c} ; z \right]. \end{aligned} \)

To verify \( (3.7) \) we substitute \( t \leftrightarrow \overline{t} \) to the integral \( (3.1) \).

**Proposition 3.2.** a) (Gauss identity) Let \( |c| - |a| - |b| > 0 \). Then

\( (3.9) \quad \begin{aligned} 2F_1^c \left[ \frac{a, b}{c} ; 1 \right] := \lim_{z \to 1} 2F_1^c \left[ \frac{a, b}{c} ; z \right] &= \frac{\Gamma^c(c) \Gamma^c(c - a - b)}{\Gamma^c(c - a) \Gamma^c(c - b)}. \end{aligned} \)

b) Let \( |c| < 1 \). Then

\( \begin{aligned} 2F_1^c \left[ \frac{a, b; c}{0} \right] := \lim_{z \to 0} 2F_1^c \left[ \frac{a, b; c}{z} \right] = 1. \end{aligned} \)

**Proof.** a) We substitute \( z = 1 \) to \( (3.1) \) and come to a beta function,

\( \pi B^c[c, c - a] / \pi B^c[b, c - a] \).

However, this argument is valid only if the beta integral \( B^c[c, c - a] \) converges.

The general statement follows from Theorem 3.9.b proved below.

b) also is reduced to a beta-function with the same problem with the domain of convergence. The general statement follows from Theorem 3.9.a. \( \square \)

**Proposition 3.3.**

\( (3.10) \quad \begin{aligned} 2F_1^c \left[ \frac{a, b; c}{z} \right] &= 2F_1^c \left[ \frac{b, a; c}{z} \right]. \end{aligned} \)

This will become obvious after Theorem 3.9. We use this symmetry in the next two proofs.

**Proposition 3.4. (Euler and Pfaff transformations),**

\( (3.11) \quad \begin{aligned} 2F_1^c \left[ \frac{a, b}{c} ; z \right] &= (1 - z)^{-a} 2F_1^c \left[ \frac{a, c - b}{c} ; \frac{z}{z - 1} \right] \end{aligned} \)

\( (3.12) \quad \begin{aligned} &= (1 - z)^{-b} 2F_1^c \left[ \frac{c - a, b}{c} ; \frac{z}{z - 1} \right] \end{aligned} \)

\( (3.13) \quad \begin{aligned} &= (1 - z)^{-c - a - b} 2F_1^c \left[ \frac{c - a, c - b}{c} ; z \right] \end{aligned} \)

**Proof.** We substitute \( t = 1 - s \) to \( (3.1) \) and get \( (3.11) \). Applying \( (3.10) \) we get \( (3.12) \). Applying \( (3.11) \) and \( (3.12) \), we get \( (3.13) \). \( \square \)

\( ^{18} \text{If } |c| \geq 1, \text{ then } \lim_{z \to 0} |2F_1^c[a, b; c; z]| = \infty. \)
Proposition 3.5. (Kummer symmetries) The following functions \( u_j^C(z) \) are equal\(^\text{19}\):

\[
(3.14) \quad u_j^C(z) = _2F_1^C \left[ \begin{array}{c} a, b \\ c \end{array} ; z \right] \tag{3.14} \quad (\text{compare with } [9], (2.2.9.1));
\]

\[
(3.15) \quad u_5^C(z) = \frac{(-1)^{c-a-b} \Gamma^C(c) \Gamma^C(c-1)}{\Gamma^C(a) \Gamma^C(b) \Gamma^C(c-a) \Gamma^C(c-b)} \zeta^{1-c} \left[ \begin{array}{c} b - c + 1; a - c + 1 \\ 2 - c \end{array} ; z \right] \tag{3.15}
\]

(see [9], (2.2.9.17)) and ratio of coefficients at \( u_1, u_5 \) in (2.2.10.35));

\[
(3.16) \quad u_3^C(z) = \frac{\Gamma^C(c) \Gamma^C(b-a)}{\Gamma^C(b) \Gamma^C(c-a)} (-z)^{-a} _2F_1^C \left[ \begin{array}{c} a, a-c+1 \\ a-b+1 ; z^{-1} \end{array} \right] \tag{3.16}
\]

(see [9], (2.2.9.9) and \( B_1 \) in (2.2.10.5));

\[
(3.17) \quad u_4^C(z) = \frac{\Gamma^C(c) \Gamma^C(a-b)}{\Gamma^C(a) \Gamma^C(c-b)} (-z)^{-b} _2F_1^C \left[ \begin{array}{c} b, b-c+1 \\ b-a+1 ; z^{-1} \end{array} \right] \tag{3.17}
\]

(see [9], (2.2.9.10) and \( B_2 \) in (2.2.10.5));

\[
(3.18) \quad u_2^C(z) = \frac{\Gamma^C(c) \Gamma^C(c-a-b)}{\Gamma^C(a-c) \Gamma^C(c-b)} _2F_1^C \left[ \begin{array}{c} a, b \\ a+b+1-c ; 1-z \end{array} \right] \tag{3.18}
\]

(see [9], (2.2.9.5) and \( A_1 \) in (2.2.10.5));

\[
(3.19) \quad u_6^C(z) = \frac{\Gamma^C(c) \Gamma^C(a+b-c)}{\Gamma^C(a) \Gamma^C(b)} (1-z)^{-c-a-b} _2F_1^C \left[ \begin{array}{c} c-a, c-b \\ c-a-b+1 ; 1-z \end{array} ; z^{-1} \right] \tag{3.19}
\]

(see [9], (2.2.9.5) and \( A_1 \) in (2.2.10.5)).

**Remark.** For each expression (3.14)–(3.19) we can apply one of the transformations (3.11)–(3.13). In this way we get 24 expressions of this type. \( \square \)

**Proof.** The formula for \( u_3 \). Changing a variable \( t = 1/s \) in (3.1) we come to

\[
\left( -1 \right)^{c-a-b-1} \zeta^{1-a} \int \frac{s^{a-c} \left( 1-s \right)^{c-b-1}}{\pi B^C(b-c-b)} d\zeta = \left( -1 \right)^{c-a-b-1} \pi B^C(a-c+1, c-b) \zeta^{-a} _2F_1^C \left[ \begin{array}{c} a, a-c+1 \\ a-b+1 ; z^{-1} \end{array} ; z^{-1} \right].
\]

We cancel \( \Gamma^C(c-b) \) and apply (2.2.3) two times.

The formula for \( u_4 \). We transpose \( a \) and \( b \) in the formula for \( u_3 \).

The formula for \( u_5 \). We combine the transformations (3.10) and (3.17).

The formula for \( u_2 \). We combine the transformations (3.10), (3.11), and again (3.10).

We combine transformations (3.10), (3.12), and again (3.10). \( \square \)

**Remark.** Proposition (3.5) is a self-closed collection of identities. However, they are reflections of the Kummer table of solutions of the hypergeometric equation

\[
\left( z(1-z) \frac{\partial^2}{\partial z^2} + (c-(a+b+1)z) \frac{\partial}{\partial z} - ab \right) u(z) = 0,
\]

\(^{19}\)The meaning of subscripts \( j \) in \( u_j^C \), references, and comments are explained in a remark after the proof.
see Erdélyi, et al., [9], Section 2.2.9, formulas (1)–(24). The Kummer table contains 6 solutions, each of them is defined in a neighborhood of one of the singular points 0, 1, ∞.

\[
\begin{align*}
&u_1(z) = \alpha_1(z), & &u_5(z) = z^{1-c}\alpha_5(z), \\
&u_3(z) = (-z)^{-a}\alpha_3(z^{-1}), & &u_4(z) = (-z)^{-b}\alpha_4(z^{-1}), \\
&u_2(z) = \alpha_2(1-z), & &u_6(z) = (1-z)^{c-a-b}\alpha_6(1-z),
\end{align*}
\]

where \(\alpha_j(x)\) are power series, \(\alpha_j(0) = 1\). Generally, these solutions are ramified at the points 0, 1, ∞. Each solution is represented in 4 forms, which can be obtained one from another by the Pfaff transformations, see Erdélyi, et al., [9], Sect. 2.1, (22)–(23). In the table above we present the corresponding expressions for \(2\mathcal{F}_1^c[a, b; c; z]\), they correspond to Kummer’s expressions with change \((a, b, c) \rightarrow (a', b', c')\). The resulting functions \(u_7^c\) are non-ramified (by definition) and differ by factors independent of \(z\), we normalize them to make them equal one to another. Counterparts of these factors (except one formula) are present in the Kummer formulas as coefficients of transfer-matrices \((u_1, u_5)\) to \((u_3, u_4)\) and \((u_2, u_6)\), with the same replacement \((a, b, c) \rightarrow (a, b, c)\), see Erdélyi, et al., [9], display (2.2.10.5) and the coefficients \(A_1, A_2, B_1, B_2\). So, in each line of Proposition 3.6 we give a reference to the corresponding formula in Erdélyi, et al., [9], (2.2.9.1)–(2.2.9.24) and to the corresponding coefficient in [9], display (2.2.10.5).

3.3. Differential equations.

Lemma 3.6.

\[
\begin{align*}
\frac{\partial}{\partial z} 2\mathcal{F}_1^c[a|a'; b|b'; c|c'; z] &= \frac{ab}{c} 2\mathcal{F}_1^c[a + 1|a'; b + 1|b'; c + 1|c'; z]; \\
\frac{\partial}{\partial z} 2\mathcal{F}_1^c[a|a'; b|b'; c|c'; z] &= \frac{ab'}{c} 2\mathcal{F}_1^c[a|a' + 1; b|b' + 1; c|c' + 1; z].
\end{align*}
\]

\[\text{Proof.}\] We differentiate the integral with respect to the parameter \(z\), and get an integral of the same form. The calculation is valid if \(\Xi \cap (\Xi + \frac{1}{2}, \frac{1}{2}) \neq \emptyset\), where \(\Xi\) is the simplex defined by (3.2)–(3.3). This intersection is open and nonempty. It remains to refer to the meromorphic continuation. \(\square\)

Denote

\[
\begin{align*}
D &= D[a, b, c] := z(1-z)\frac{\partial^2}{\partial z^2} + (c - (a + b + 1)z)\frac{\partial}{\partial z} - ab; \\
D' &= D'[a', b', c'] := (1-\Xi)\frac{\partial^2}{\partial \xi^2} + (c' - (a' + b' + 1)\Xi)\frac{\partial}{\partial \xi} - a'b'.
\end{align*}
\]

Proposition 3.7. The complex hypergeometric function \(\mathcal{F}(z) = 2\mathcal{F}_1^c[a, b; c; z]\) satisfies the following system of partial differential equations

\[
\begin{align*}
D[a, b, c]\mathcal{F} &= 0; & D'[a', b', c']\mathcal{F} &= 0.
\end{align*}
\]

We call these equation by complex hypergeometric system.

\[\text{Proof.}\] This follows from the identity

\[
D[a, b, c]\left(t^{b-1}(1-t)^{a-b-1} - t z^{a-1}\right) = -a\frac{\partial}{\partial t}\left(t^b(1-t)^{c-b} - t z^a\right)
\]
(cf. [9], (2.1.3.11)). Consider sufficiently small positive \( \varepsilon, \delta, \kappa \) and take \( a, b, c \) such that

\[
[b]' = \varepsilon, \quad [c]' = \varepsilon + \delta, \quad [a]' = -\frac{1}{2} + \varepsilon + \delta + \kappa.
\]

We multiply both parts of (3.23) by \( (1 - \tau)^{c - b - 1}(1 - \overline{\tau})^{-c'} \) and integrate over \( \mathbb{C} \). In the left hand side for such values of the parameter we can permute integration and differentiation in \( z \). In the right hand side the integrand is an integrable derivative of an integrable function. Therefore the right hand side is zero.

**Proposition 3.8.** a) Any solution of system (3.20)–(3.21) is real analytic in \( z \).

b) Let \( z_0 \neq 0, 1, \infty \). Denote by \( \varphi_1(z), \varphi_2(z) \) a pair of independent holomorphic solutions of the ordinary differential equation \( D[a, b, c] f(z) = 0 \) at a neighborhood of \( z_0 \). Denote by \( \psi_1(\tau), \psi_2(\overline{\tau}) \) a pair of antiholomorphic solutions of the ordinary differential equation \( D'[a', b', c'] f(\tau) = 0 \). Then any solution of the system (3.20)–(3.21) can be represented as

\[
\sum_{i,j=1,2} \tau_{ij}(a, b, c) \varphi_i(z) \psi_j(\tau).
\]

c) If we choose \( \varphi_i, \psi_j \) meromorphic in the parameters \( a, b, c \) in some domain in \( \mathbb{C}^3 \), then the coefficients \( \tau_{ij} \) also are meromorphic in the parameters \( a, b, c \).

**Proof.** a) Indeed, \( D[a, b, c] \) is an elliptic differential operator, therefore solutions of the equation \( Df = 0 \) are analytic functions, see, e.g., [19], Theorem 9.5.1.

b) Consider a solution

\[
\sum_{i,j=1,2} \tau_{ij}(a, b, c) \varphi_i(z) \psi_j(\tau).
\]

of the system of partial differential equations (3.22). These equations determine recurrence relations for the Taylor coefficients \( h_{ij} \) of \( \frac{F(z)}{F(z_0)} \) at \( z_0 \). It can easily be checked that all the coefficients \( h_{ij} \) admit linear expressions in terms of \( h_{00}, h_{01}, h_{10}, h_{11} \). On the other hand, for given \( h_{00}, h_{01}, h_{10}, h_{11} \), we can find a local solution of the complex hypergeometric system (3.22) in the form \( \sum C_{ij} \varphi_i(z) \psi_j(\tau) \).

c) By Lemma 3.7, the coefficients \( h_{00}, h_{10}, h_{01}, h_{11} \) depend on \( a, b, c \) meromorphically. If \( \varphi_i(z_0), \varphi'_i(z_0), \psi_j(z_0), \psi'_j(z_0) \) are meromorphic in the parameters, then the coefficients \( C_{ij} \) also are meromorphic.

**3.4. Expressions for \( 2F_1F \).** Let us write expansions of \( 2F_1F[\ldots; z] \) near the singular points \( z = 0, 1, \infty \). Explicit formulas for fundamental systems of solutions of the hypergeometric differential equation are well-known, see Erdélyi, et al., [9], 2.9 (the Kummer series). For definiteness, consider \( z_0 = 0 \). If \( c \notin \mathbb{Z} \), then for generic values of the parameters the hypergeometric equation \( D[a, b, c] f(z) = 0 \) has two holomorphic (ramified) solutions on a punctured neighborhood of 0,

\[
\varphi_1(z) = 2F[a + 1 - c, b + 1 - c; z], \quad \varphi_2(z) = z^{1-c} F[a + 1 - c, b + 1 - c; 2 - c; z].
\]

The equation \( D'[a', b', c'] f(\tau) = 0 \) has two antiholomorphic solutions

\[
\psi_1(z) = 2F[a', b', c', \tau], \quad \psi_2(z) = \tau^{1-c'} F[a' + 1 - c', b' + 1 - c'; 2 - c'; \tau].
\]

Therefore near \( z = 0 \) we have solutions of system (3.20)–(3.21) of the same form (3.24) with new \( \varphi, \psi \). We get a family of functions depending of 4 parameters \( \tau_{ij} \), therefore for generic \( a, b, c \) this formula gives all multivalued solutions near \( z = 0 \).
Solutions (3.24) that are single valued in a neighborhood of 0 have the form
\[ (3.25) \quad A_1 \varphi_1(z) \psi_1(\overline{z}) + A_2 \varphi_2(z) \psi_2(\overline{z}). \]

**Theorem 3.9.** a) In the disk $|z| < 1$ we have the following expansion:

\[ (3.26) \quad 2F_1^C \left[ \begin{array}{c} a, b \\ c \end{array} ; z \right] = A_0 \cdot 2F_1 \left[ \begin{array}{c} a, b \\ c \end{array} ; z \right] 2F_1 \left[ \begin{array}{c} a', b' \\ c' \end{array} ; z \right] + \\
+ A_1 \cdot z^{1-c|1-c'} 2F_1 \left[ \begin{array}{c} a + 1 - c, b + 1 - c \\ 2 - c \end{array} ; z \right] 2F_1 \left[ \begin{array}{c} a' + 1 - c', b' + 1 - c' \\ 2 - c' \end{array} ; z \right], \]

where
\[ (3.27) \quad A_0 = 1, \]
\[ (3.28) \quad A_1 = (-1)^{c-a-b} \frac{\Gamma^C(c) \Gamma^C(c-1)}{\Gamma^C(a) \Gamma^C(b) \Gamma^C(c-a) \Gamma^C(c-b)}. \]

b) In the disk $|z-1| < 1$ the following expansion holds:

\[ (3.29) \quad 2F_1^C \left[ \begin{array}{c} a, b \\ c \end{array} ; z \right] = B_0 \cdot 2F_1 \left[ \begin{array}{c} a, b \\ a+b+1-c;1-z \end{array} ; z \right] 2F_1 \left[ \begin{array}{c} a', b' \\ a'+b'+1-c';1-z \end{array} ; z \right] + \\
+ B_1 \cdot (1-z)^{c-a-b|c-a-b'} 2F_1 \left[ \begin{array}{c} c - a, c - b \\ c+1-a-b;1-z \end{array} ; z \right] 2F_1 \left[ \begin{array}{c} c' - a', c' - b' \\ c'+1-a'-b';1-z \end{array} ; z \right], \]

where
\[ (3.30) \quad B_0 = \frac{\Gamma^C(c) \Gamma^C(c-a-b)}{\Gamma^C(c-a) \Gamma^C(c-b)}, \]
\[ (3.31) \quad B_1 = \frac{\Gamma^C(c) \Gamma^C(a+b-c)}{\Gamma^C(a) \Gamma^C(b)}. \]

c) In the disk $|z| > 1$ the following expansion holds:

\[ (3.32) \quad 2F_1^C \left[ \begin{array}{c} a, b \\ c \end{array} ; z \right] = C_0 \cdot (-z)^{-a|c-a'} 2F_1 \left[ \begin{array}{c} a, a+1-c \\ a+1-b;1-z \end{array} ; z \right] 2F_1 \left[ \begin{array}{c} a', a'+1-c' \\ a'+1-b';1-z \end{array} ; z \right] + \\
+ C_1 \cdot (-z)^{-b|c-b'} 2F_1 \left[ \begin{array}{c} b, b+1-c \\ b+1-a;1-z \end{array} ; z \right] 2F_1 \left[ \begin{array}{c} b', b'+1-c' \\ b'+1-a';1-z \end{array} ; z \right], \]

where
\[ (3.33) \quad C_0 := \frac{\Gamma^C(c) \Gamma^C(b-a)}{\Gamma^C(c-a) \Gamma^C(b)}, \]
\[ (3.34) \quad C_1 := \frac{\Gamma^C(c) \Gamma^C(a-b)}{\Gamma^C(c-b) \Gamma^C(a)}. \]

### 3.5. Proof of Theorem 3.9

Forms (3.24), (3.28), (3.31) for the desired expressions follow from the preceding considerations. Also we know that the coefficients $A_0$, $A_1$, $B_0$, $B_1$, $C_0$, $C_1$ are meromorphic in $a$, $b$, $c$. Now we apply asymptotic expansions from Theorem 2.3.

1. Asymptotic of $2F_1^C[a, b; c; z]$ as $z \to \infty$. Assume that the defining integral for $2F_1^C[a, b; c; z]$ converges, and also
\[ (3.35) \quad |b - a| > 0 \]
Then
\[
\frac{1}{\pi B^C(b, c - b)} \int_C t^{b-1}(1-t)^{c-b-1}(1-zt)^{-a} d\bar{t} =
\]
\[
= \frac{z^{-a}}{\pi B^C(b, c - b)} \int_C t^{b-1}(1-t)^{c-b-1}(z^{-1} - t)^{-a} d\bar{t} \sim
\]
\[
\frac{B^C(b - a, c - b)}{B^C(b, c - b)} \cdot (-z)^{-a} \cdot \left(1 + \sum_{(i,i') \neq (0,0)} p_{ii'} z^{-i-i'} \right) +
\]
\[
+ \frac{B^C(b, 1-a)}{B^C(b, c - b)} \cdot z^{-b} \cdot \left(1 + \sum_{(i,i') \neq (0,0)} q_{ii'} z^{-i-i'} \right).
\]

Precisely, denote \(z^{-1}\) by \(\varepsilon\), and denote the integrand in the last integral by \(H(\cdot)\).

Let \(\varphi(t) \geq 0, \psi(t) \geq 0\) be smooth functions such that \(\varphi(t) + \psi(t) = 1\), \(\varphi(t) = 1\) near 0, and \(\psi(t) = 1\) near \(\infty\). A straightforward differentiation with respect to the parameter \(\varepsilon\) shows that
\[
\int H(t; \varepsilon) \psi(t) d\bar{t}
\]
is smooth near \(\varepsilon = 0\). For
\[
\int H(t; z) \varphi(t; \varepsilon) d\bar{t}
\]
we apply Theorem 2.13 due to the restriction \[(3.35)\] we can also apply (2.14). Thus
we get explicit coefficients \(C_0, C_1\) in the expansion \[(3.32)\] To remove restrictions
for the parameters, we refer to the analytic continuation.

Finally, we transform \(B^C(b, 1-a)\) with formula \[(2.3)\],
\[
B^C(b, 1-a) = \frac{\Gamma^C(b)}{\Gamma^C(1+b-a)} = \frac{(-1)^b \Gamma^C(b) \Gamma^C(a-b)}{\Gamma^C(a)}.
\]

2. Asymptotic of \(2F_1^C[a, b; c; z]\) as \(z \to 0\). Substituting \(t = 1/s\) to the definition (1.18) of \(2F_1^C\), we get
\[
2F_1^C[a, b; c; z] = \frac{(-1)^{c-a-b}}{\pi B^C(b, c - b)} \int_C s^{-c+a}(z-s)^{-a}(1-s)^{c-b-1} d\bar{s} \sim
\]
\[
\sim \frac{(-1)^{c-a-b} B^C(a - c + 1, 1 - a)}{B^C(b, c - b)} \cdot z^{1-c} \cdot \left(1 + \sum_{(i,i') \neq (0,0)} p_{ii'} z^{ij} \right) +
\]
\[
+ \frac{(-1)^{c-b} B^C(1 - c, c - b)}{B^C(b, c - b)} \cdot \left(1 + \sum_{(i,i') \neq (0,0)} q_{ii'} z^{ij} \right).
\]
3. Asymptotic of $\, _2F_1^C[a,b;c;z]$ as $z \to 1$. We substitute $t = \frac{1}{1-z}$ to (1.18) and get

$$\, _2F_1^C[a,b;c;z] = \frac{(-1)^{c-b}}{\pi B^C(b,c-b)} \int_C s^{c-b-1}(1-s)^{a-c}(1-z-s)^{-a} \, ds \sim$$

$$\sim \frac{(-1)^{c-b}B^C(c-b,1-a)}{B^C(b,c-b)} \cdot (1-z)^{c-b-a} \cdot \left(1 + \sum_{(i,j) \neq (0,0)} p_{ij}(1-z)^{|ij|} + \frac{(-1)^{c-b-a}B^C(c-b,1+a-c)}{B^C(b,c-b)} \cdot \left(1 + \sum_{(j,j') \neq (0,0)} q_{ij}z^{|j|}\right)\right) .$$

**Remark. Another way of a proof of Theorem 3.9.** Applying the Kummer formulas, Erdélyi, et al., [9], Section 2.9, we can write the analytic continuation of (3.25) to a neighborhood of this point. The resulting expression for $\, _2F_1^C$ must be non-ramified at $z = 1$. This gives us the coefficients in (3.25) up to a common factor. In fact this calculation is done below in the proof of Proposition 3.11. The scalar factor can be evaluated using (3.37). It remains to apply the Kummer formulas ([9], Section 2.9) for the analytic continuation again and to get an expansion at $\infty$. □

3.6. Additional symmetry.

**Proposition 3.10.** Let $a - b \in \mathbb{Z}$. Then

$$\, _2F_1^C \left[ \begin{array}{c} a \cr b \end{array} \right] c', z \right] = \, _2F_1^C \left[ \begin{array}{c} a' \cr b' \end{array} \right] c', z \right] .$$

**Proof.** The expansions (3.26)–(3.28) at 0 for both functions are identical. We only must verify the equality of the denominators in (3.28):

$$\Gamma^C(a|a') \Gamma^C(b|b') \Gamma^C(c-a|c-a') \Gamma^C(c-b|c-b') = \Gamma^C(b|b') \Gamma^C(c-a|c-a') \Gamma^C(c-b|c-b') \Gamma^C(c-b|c-b') .$$

The both sides are equal to

$$\frac{(-1)^{c'-c} \pi^4 \Gamma(a) \Gamma(a') \Gamma(b) \Gamma(b') \Gamma(c-a) \Gamma(c-a') \Gamma(c-b) \Gamma(c-b')}{\sin \pi a' \sin \pi b' \sin \pi (c'-a') \sin \pi (c'-b')} .$$

□

3.7. Degenerations and logarithmic expressions.

a) Residues and zeros. Notice that poles and zeros of $\, _2F_1^C[a,b;c;z]$ as a function of $a, b, c$ depend on a choice of a normalizing factor in the front of the integral (3.1).

It is easy to see that residues at poles also are solutions of the complex hypergeometric system (3.22). The expressions for the residues can be obtained from our expansions.

For obtaining the residues at $\{a|a'\} \in \mathbb{N} \times \mathbb{N}$ we can use the expansion of $\, _2F_1^C$ at $z = 0$, see (3.26)–(3.28). We get

$$z^{1-c-1-c'} \, _2F_1 \left[ \begin{array}{c} a + 1 - c, b + 1 - c' \cr 2 - c \end{array} \right] z \, _2F_1 \left[ \begin{array}{c} a' + 1 - c', b' + 1 - c' \cr 2 - c' \end{array} \right]$$

with an obvious $\Gamma^C$-factor. Applying the Pfaff transformations of $\, _2F_1$, we observe that these expressions are elementary functions. Formulas (3.26)–(3.28) allow to calculate residues at the poles of all the types (3.4).
Next, consider another normalization\(^{20}\) of the functions \(\tilde{2F}_1^C\):

\[
\tilde{F}_1^C[a, b; c; z] := \frac{1}{\Gamma_C(c)} 2F_1^C[a, b; c; z].
\]

This operation cancels the factor \(\Gamma_C(c)\) in expansion of \(2F_1^C[z]\) at \(\infty\), see (3.32)–(3.34). So we get a finite expression at the poles (3.5) and non-zero function at the zeros (3.6).

Thus, at all exceptional planes (3.4)–(3.6) we get explicit nonzero expressions. Such expressions also depend on normalization of \(2F_1^C[...; z]\) but for a point \((a_0, b_0, c_0)\) being in a general position on an exceptional plane such nonzero expression is canonically defined up to a constant factor.

b) Further degenerations. Classical hypergeometric differential equation has a sophisticated list of degenerations, see [9], Sect. 2.2. In our case new difficulties arise if at least two of the parameters \(a, b, c - a, c - b\) are contained in \(\mathbb{Z} \times \mathbb{Z}\). We stop here further analysis and only notice that for exceptional values of the parameters a solution of the complex hypergeometric system (3.22) can be non-unique.

For instance, if \(a \in \mathbb{Z}_- \times \mathbb{Z}_-\), \(c - b \in \mathbb{N} \times \mathbb{N}\), then both summands in (3.25) are single-valued (since all hypergeometric series are terminating).

c) Logarithmic expressions. For definiteness we discuss the case \(c \in \mathbb{N} \times \mathbb{N}\) (which is interesting for our further purposes). Consider the function \(\tilde{2F}_1^C\) defined by (3.38). It has a removable singularity at our \(c\). Recall that for \(c = n \in \mathbb{N}\) the usual hypergeometric differential equation \(D[a, b, n] f = 0\) has two solutions. The first is \(2F_1[a, b; n; z]\) and the second has the form

\[
\Psi[a, b; n; z] = \sum_{j=-n+1}^{\infty} p_j z^j + \ln z \cdot 2F_1[a, b; n; z],
\]

where \(p_j\) are explicit coefficients, \(p_{-n+1} \neq 0\), and this form does not depend on further degenerations, see [1], Section 2.3. Passing around 0 we get

\[
\Psi[a, b; n; e^{i\varphi} z] \bigg|_{\varphi = 2\pi} = \Psi[a, b; n; z] = 2F_1[a, b; n; z].
\]

Thus the system

\[
D[a, b; n] \mathcal{F} = 0, \quad D[a', b'; n'] \mathcal{F} = 0
\]

has two solutions that are single-valued near zero, the first is obvious

\[2F_1[a, b, n; z] 2F_1[a', b'; n'; \mathcal{F}],\]

and the second is

\[
2F_1[a, b; n; z] \Psi[a', b'; n'; \mathcal{F}] + \Psi[a, b; n; z] 2F_1[a', b'; n'; \mathcal{F}].
\]

Our function \(\tilde{2F}_1^C[a, b; n; z]\) is certain linear combination of these solutions.

d) On uniqueness of a solution of the hypergeometric system.

---

\(^{20}\)In fact, in the main part of our work we use this normalization of the kernel, see (1.19). Due to this we do not lose the case of \(L^2\) on the complex quadric discussed in Subsect. 1.3.
Proposition 3.11. Let
\[a, b, c, c - a - b, c - a, c - b \notin \mathbb{Z},\]
\[a', b', c', c' - a' - b', c' - a', c' - b' \notin \mathbb{Z}.
\]
Let the system \(D[a, b, c]f = 0, D'[a', b', c']f = 0\) have a non-ramified non-zero solution. Then \(c - c' \in \mathbb{Z}\) and
\[a - a', b - b' \in \mathbb{Z} \text{ or } a - b', b - a' \in \mathbb{Z}.
\]
Such solution is unique up to a scalar factor and therefore is \(2F_1^c[a|b|b'; c'|c'; z] \text{ or } 2F_1^c[a|b'|b'; c|c'; z] \)

Proof. First, we examine the behavior of a solution near \(z = 0\). Let
\[
\varphi(z) := 2F_1 \left[ \frac{a, b}{c} : z \right], \quad \psi(z) := 2F_1 \left[ \frac{a + 1 - c, b + 1 - c}{2 - c} : z \right],
\]
i.e., \(\varphi, z^{1-c}\psi\) are the Kummer solutions of the hypergeometric equation \(D[a, b, c]f = 0\) at 0, see [9], (2.9.1), (2.9.17). Denote by \(\varphi(\tau), \psi(\tau)\) the similar functions obtained by the change \(a \mapsto a', b \mapsto b', c \mapsto c', z \mapsto \tau\). A solution of our system near 0 has the form
\[
G(z) = \sigma \varphi(z)\varphi(\tau) + \mu z^{1-c} \varphi(z)\psi(\tau) + \nu z^{1-c} \psi(z)\varphi(\tau) + \tau z^{1-c} \psi(z)\psi(\tau).
\]
Passing \(m\) times around 0 we come to
\[
G(z) = \sigma \varphi(z)\varphi(\tau) + \mu e^{2\pi mc'i} z^{1-c} \varphi(z)\psi(\tau) + \nu e^{2\pi mc'i} z^{1-c} \psi(z)\varphi(\tau) + \tau e^{2\pi mc'i} z^{1-c} \psi(z)\psi(\tau).
\]
Since \(c, c' \notin \mathbb{Z}\), we have \(e^{2\pi mc'i}, e^{-2\pi mc'i} \neq 1\), on the other hand they are \(\neq e^{2\pi mc'i}\). If \(G(z)\) is single-valued, then \(\mu = \nu = 0\). Also, we need \(\tau = 0\) or \(c - c' \in \mathbb{Z}\).

To examine the behavior of \(G\) near \(z = 1\) we apply a formula for analytic continuation, see [9], Subsect. 2.10. Near \(z = 1\) we have
\[
F \left[ \frac{a, b}{c} : z \right] = A_1(a, b, c) 2F_1 \left[ \frac{a, b}{a + b - c + 1} : 1 - z \right] + \nonumber
\]
\[
+ A_2(a, b, c) (1 - z)^{c-a-b} 2F_1 \left[ \frac{c - a, c - b}{c - a - b + 1} : 1 - z \right],
\]
where
\[
A_1(a, b, c) := \frac{\Gamma(c)\Gamma(c - a - b)}{\Gamma(c - a)\Gamma(c - b)}, \quad A_2(a, b, c) := \frac{\Gamma(c)\Gamma(a + b - c)}{\Gamma(a)\Gamma(b)}.
\]
Since \(c - a - b, c' - a' - b' \notin \mathbb{Z}\), the expression \(\varphi(z)\varphi(\tau)\) is not single valued. Thus \(\tau \neq 0\), \(c - c' \in \mathbb{Z}\), and
\[
G(z) = \sigma \varphi(z)\varphi(\tau) + \tau z^{1-c} \psi(z)\psi(\tau).
\]
Applying for \(\varphi, \varphi, \psi, \psi\) formula (3.42) and the identity
\[
2F_1(\alpha, \beta; \gamma; z) = (1 - z)^{\gamma - a - \beta} 2F_1(\gamma - \alpha, \gamma - \beta; \gamma; z),
\]
we come to

\[ G(z) = \sigma A(a, b, c) A(a', b', c') {\, _2F_1} \left[ \begin{array}{c} a, b \\ c \\ \end{array} ; 1 - z \right] {\, _2F_1} \left[ \begin{array}{c} a', b' \\ c' \\ \end{array} ; 1 - \frac{1}{c} \right] + \]

\[ + \left\{ \sigma A_1(a, b, c) A_2(a', b', c') + \tau A_1(a+1-c, b+1-c, 2-c) A_2(a'+1-c', b'+1-c', 2-c') \right\} \times \]

\[ \times (1 - \frac{1}{c})^{a-a'-b'} {\, _2F_1} \left[ \begin{array}{c} c-a, c-b \\ c-a-b+1 \\ \end{array} ; 1 - z \right] {\, _2F_1} \left[ \begin{array}{c} c', b' \\ c' \\ \end{array} ; 1 - \frac{1}{c} \right] + \]

\[ + \left\{ \sigma A_2(a, b, c) A_1(a', b', c') + \tau A_2(a+1-c, b+1-c, 2-c) A_1(a'+1-c', b'+1-c', 2-c') \right\} \times \]

\[ \times (1 - z)^{c-a-b} {\, _2F_1} \left[ \begin{array}{c} c-a, c-b \\ c-a-b+1 \\ \end{array} ; 1 - z \right] {\, _2F_1} \left[ \begin{array}{c} c', b' \\ c' \\ \end{array} ; 1 - \frac{1}{c} \right] + \]

\[ + A_2(a+1-c, b+1-c, 2-c) A_1(a'+1-c', b'+1-c', 2-c') \times \]

\[ \times (1 - z)^{c-a-b'} {\, _2F_1} \left[ \begin{array}{c} c-a, c-b \\ c-a-b+1 \\ \end{array} ; 1 - z \right] {\, _2F_1} \left[ \begin{array}{c} c', b' \\ c' \\ \end{array} ; 1 - \frac{1}{c} \right]. \]

The coefficients \( A_1(\cdot), A_2(\cdot) \) have no zeros and no poles under our restrictions. The expression is single-valued if and only if two curly brackets are zero and \( (c - a - b) - (c' - a' - b') \in \mathbb{Z} \). This implies

\[ (a + b) - (a' + b') \in \mathbb{Z}. \]

Two curly brackets give a system of linear equations for \( \sigma, \tau \). It has a nonzero solution if and only if its determinant \( \Delta \) is zero. Straightforward calculations give

\[ \Delta = \pi^{-3} \Gamma(c) \Gamma(c') \Gamma(2-c) \Gamma(2-c') \Gamma(c-a-b) \Gamma(c-a'-b') \Gamma(a+b-c) \Gamma(a'+b'-c') \Xi, \]

where

\[ \Xi = \sin \pi(c-a) \sin \pi(c-b) \sin \pi a' \sin \pi b' - \sin \pi(c'-a') \sin \pi(c'-b') \sin \pi a \sin \pi b. \]

Clearly, the set \( \Xi(a, b, c, a', b', c') = 0 \) is invariant with respect to the shifts \( a \mapsto a+1, b \mapsto b+1, c \mapsto c+1 \). Therefore to examine the set of zeros we can assume \( c' = c, b' = a + b - a' \). Under these conditions \( \Xi \) can be reduced to the following form:

\[ \Xi(a, b, c, a', b', c') = \sin \pi(a-a') \sin \pi(a'-b) \sin \pi c \sin \pi(c-a-b) \]

(this non-obvious identity can be verified by decompositions of both sides into sums if exponentials). This implies \( \ref{3.31} \).

If \( \Delta = 0 \) then \( \sigma, \tau \) are defined up to a common scalar factor, this proves the uniqueness (and gives an expression for \( \sigma/\tau \)).

E) NON-INTERESTING SOLUTIONS. However, we have seen that the complex hypergeometric system for some values of the parameters has two single-valued solutions. Also, there are solutions that do not seem reasonable. For instance, we have

\[ D[0, b_1, c_1] \cdot 1 = 0, \quad D'[0, b_2, c_2] \cdot 1 = 0 \]

for arbitrary \( b_1, c_1, b_2, c_2 \in \mathbb{C} \).

3.8. Differential-difference equations for \( \, _2F_1^{\mathbb{C}} \). We can regard \( A_F[1] = \, _2F_1[a, b; c; z] \) as a family of functions of a complex variable \( z \) depending on 3 parameters \( a, b, c \). But we also can regard \( A_F[2] = \, _2F_1[a, b; c; z] \) as one function of the four complex variables \( a, b, c, z \). Then \( \, _2F_1[a, b; c; z] \) satisfy a non-obvious system of linear differential-difference equations, some examples of such equations are in Erdélyi, et al., \[39\], (2.8.20-45). Below we show that such equations can be automatically transformed
to differential-difference equations for the function \( _2F_1^C[a|a', b|b'; c|c'; z] \) of 7 complex variables.

Consider a space of functions in the variables \( a, b, c, z \). Define operators

\[
T_a f(a, b, c, z) = f(a + 1, b, c, z), \quad T_b f(a, b, c, z) = f(a, b + 1, c, z), \quad T_c f(a, b, c, z) = f(a, b, c + 1, z).
\]

Consider finite sums of the form

\[
(3.44) \quad \mathcal{L} = \sum_{j \geq 0} \sum_{k,l,m} U_{j,k,l,m}(a, b, c, z) T_a^j T_b^l T_c^m \frac{\partial^j}{\partial z^j},
\]

where \( U_{j,k,l,m}(a, b, c, z) \) are polynomial expressions in \( z \) with coefficients rationally depending on \( a, b, c \).

Assume that

\[
\mathcal{L} \ _2F_1[a, b; c; z] = 0.
\]

We can regard an operator (3.44) as an operator on functions \( f(a|a', b|b', c|c', z) \) on \( \Lambda^3 \times \hat{C} \). We also define operators

\[
T_{a'} f(a'|a', b|b', c|c', z) = f(a'|a' + 1, b|b', c|c', z),
\]

\[
T_{b'} f(a'|a', b|b', c|c', z) = f(a'|a', b|b' + 1, c|c', z),
\]

\[
T_{c'} f(a'|a', b|b', c|c', z) = f(a'|a', b|b', c|c' + 1, z).
\]

For such an operator \( \mathcal{L} \) we define the operator \( \mathcal{L}' \) by

\[
\mathcal{L}' = \sum_{j \geq 0} \sum_{k,l,m} U_{j,k,l,m}(a', b', c', z) T_a^j T_b^l T_c^m \frac{\partial^j}{\partial z^j}.
\]

From the definition it follows that

\[
\mathcal{L} \mathcal{L}' = \mathcal{L}' \mathcal{L}.
\]

**Proposition 3.12.** Let the function \( Q(a, b, c, z) = _2F_1[a, b; c; z] \) satisfy an equation \( \mathcal{L} Q = 0 \). Then the function

\[
R(a|a', b|b', c|c', z) := _2F_1^C[a|a', b|b', c|c', z]
\]

satisfies the system of equations

\[
(3.45) \quad \mathcal{L} R(a|a', b|b', c|c', z) = 0, \quad \mathcal{L}' R(a|a', b|b', c|c', z) = 0.
\]

**Lemma 3.13.** Let \( Q = _2F_1[a, b; c; z] \) satisfy an equation \( \mathcal{L} Q = 0 \). Then

\[
(3.46) \quad e^\pi c|c-b| \Gamma(c) \Gamma(c-1) \Gamma(a) \Gamma(b) \Gamma(c-a) \Gamma(c-b) z^{1-c} \ _2F_1[a+1-c, b+1-c; 2-c, z]
\]

satisfies the same equation.

**Remark.** The same statement holds for the functions

\[
(3.47) \quad u_1 = \frac{\Gamma(c-a) \Gamma(c-a-b)}{\Gamma(c) \Gamma(c-b)} _2F_1\left[\begin{array}{c}a, b \\ a+b-c+1\end{array}; 1-z\right];
\]

\[
(3.48) \quad u_2 = \frac{\Gamma(c) \Gamma(a+b-c)}{\Gamma(a) \Gamma(b)} (1-z)^{c-a-b} _2F_1\left[\begin{array}{c}c-a, c-b \\ c-a-b+1\end{array}; 1-z\right];
\]

\[
(3.49) \quad u_3 = \frac{\Gamma(c) \Gamma(b-a)}{\Gamma(b) \Gamma(c-b)} z^{-a} _2F_1\left[\begin{array}{c}a, 1-c+a \\ 1-b+a\end{array}; z^{-1}\right].
\]
and also for other summands in the right-hand sides of formulas Erdélyi, et al. [9], (2.10.1)–(2.10.4).

Proof of Lemma 3.13. First, let \( a, b, c \) be in a general position. By Erdélyi, et al. [9], (2.10.1), (2.10.5),

\[
(3.49) \quad F(a, b; c; z) = u_1 + u_2,
\]

where \( u_1, u_2 \) are given by (3.47). The function \( u_2 \) is ramified at \( z = 1 \). Passing around this point we get a function

\[
\tilde{F} := u_1 + e^{2\pi i (c-a-b)} u_2.
\]

By analytic continuation, \( \mathcal{L} \tilde{F} = 0 \). The factor \( e^{2\pi i (c-a-b)} \) does not change under the shifts \( T_a, T_b, T_c \). Therefore the summands \( u_1, u_2 \) satisfy the same equation, \( \mathcal{L} u_1 = 0, \mathcal{L} u_2 = 0 \). We apply the same transformation \( (3.49) \) to the summand \( u_1 \) and repeat the same reasoning. We observe that

\[
\frac{\pi \Gamma(c) \Gamma(c-1)}{\Gamma(c-c) \Gamma(c-b) \sin \pi(a+b-c)} z^{1-c} 2F_1[a+1-c, b+1-c; 2-c, z]
\]

satisfies the same equation. This expression differs from \( (3.49) \) by the factor \( e^{i\pi(a+b-c)} \sin \pi(a+b-c) \), which is invariant under the shifts \( T_a, T_b, T_c \).

Passing to a limit we omit restrictions to \( a, b, c \).

Proof of Proposition 3.12. We use the expression \( (3.26) \) for \( \tilde{F} \), and the expression

\[
\frac{e^{\pi i (c-a-b)} \Gamma(c) \Gamma(c-1)}{\Gamma(a) \Gamma(b) \Gamma(c-a) \Gamma(c-b)} z^{1-c} 2F_1 \left[ a+1-c, b+1-c, \frac{2-c}{c} z \right] \times \\
\frac{e^{\pi i (c'-a' - b') \Gamma(c') \Gamma(c'-1)}}{\Gamma(a') \Gamma(b') \Gamma(c'-a') \Gamma(c'-b')} \frac{z^{1-c'} 2F_1 \left[ a'+1-c', b'+1-c', \frac{2-c'}{c'-1} z \right]}{\sin \pi c \sin \pi c'}
\]

satisfies the system \( (3.44) \). It differs from the second summand in \( (3.26) \) by a factor

\[
\frac{i^0 \sin \pi a' \sin \pi b' \sin \pi (c'-a') \sin \pi (c'-b')}{\sin \pi c \sin \pi (c'-1)}
\]

This expression is invariant with respect to shifts \( T_a, T_{a'}, \ldots \). Therefore the second summand in \( (3.26) \) also satisfies the system.

3.9. One difference operator. By [34], formula (2.3), the Gauss hypergeometric function \( 2F_1[p, q; r; z] \) satisfies the following difference equation

\[
(3.50) \quad -z \cdot 2F_1(p, q; r; z) = \\
\quad = \frac{q(r-p)}{(q-p)(1+q-p)} 2F_1(p-1, q+1; r; z) - \\
\quad - \left[ \frac{q(r-p)}{(q-p)(1+q-p)} + \frac{p(r-q)}{(p-q)(1+p-q)} \right] 2F_1(p, q; r; z) + \\
\quad + \frac{p(r-q)}{(p-q)(1+p-q)} 2F_1(p+1, q-1; r; z).
\]
Define the difference operators \( L, L' \) acting on functions of the variables \( a, b, c, z \) by
\[
\begin{align*}
L &= \frac{b(c - a)}{(b - a)(1 + b - a)} (T_a^{-1}T_b - 1) - \frac{a(c - b)}{(a - b)(1 + a - b)} (T_aT_b^{-1} - 1); \\
L' &= \frac{b'(c' - a')}{(b' - a')(1 + b' - a')} (T_{a'}^{-1}T_{b'} - 1) - \frac{a'(c' - b')}{(a' - b')(1 + a' - b')} (T_{a'}T_{b'}^{-1} - 1).
\end{align*}
\]

Corollary 3.14. The complex hypergeometric function \( \text{$_2F_1^C[a, b; c; z]$} \) satisfies the following system of difference equations
\[
\begin{align*}
L \text{$_2F_1^C[a, b; c; z]$} &= z \text{$_2F_1^C[a, b; c; z]$}; \\
L' \text{$_2F_1^C[a, b; c; z]$} &= \mp \text{$_2F_1^C[a, b; c; z]$}.
\end{align*}
\]

3.10. Some properties of the kernel \( \mathcal{K} \). We have the following corollaries from our previous considerations.

1) By \( \text{(3.10)} \) \( \mathcal{K}_{a,b} \) is even,
\[
\mathcal{K}_{a,b}(z; -k, -\sigma) = \mathcal{K}_{a,b}(z; k, \sigma).
\]

2) By \( \text{(4.8)} \),
\[
\mathcal{K}_{a,b}(z; k, -\sigma) = \mathcal{K}_{a,b}(z; k, \sigma).
\]

In particular, \( \mathcal{K}_{a,b}(z; k, \sigma) \) is real on \( \Lambda \).

3) By Proposition \( \text{(3.3)} \) \( \mathcal{K}_{a,b}(z; k, \sigma) \) satisfies the following differential equations:
\[
\begin{align*}
\mathcal{D} \mathcal{K}_{a,b}(z; k, \sigma) &= \frac{1}{4}(k + \sigma)^2 \mathcal{K}_{a,b}(z; k, \sigma); \\
\overline{\mathcal{D}} \mathcal{K}_{a,b}(z; k, \sigma) &= \frac{1}{4}(k - \sigma)^2 \mathcal{K}_{a,b}(z; k, \sigma).
\end{align*}
\]

4) By Corollary \( \text{(3.14)} \) \( \mathcal{K}_{a,b}(z; k, \sigma) \) satisfies the following difference equations:
\[
\begin{align*}
\mathcal{D} \mathcal{K}_{a,b}(z; k, \sigma) &= z \mathcal{K}_{a,b}(z; k, \sigma); \\
\overline{\mathcal{D}} \mathcal{K}_{a,b}(z; k, \sigma) &= \mp \mathcal{K}_{a,b}(z; k, \sigma).
\end{align*}
\]

4. Nonexistence of commuting self-adjoint extensions

Here we prove that for \( (a, b) \notin \Pi \) the operators \( \frac{1}{2}(\mathcal{D} + \overline{\mathcal{D}}), \frac{1}{2}(\mathcal{D} - \overline{\mathcal{D}}) \) defined on \( \mathcal{D}(\hat{\mathcal{C}}) \) do not admit commuting self-adjoint extensions. We analyze the set of possible generalized eigenfunctions and show that this set is too small.

4.1. Generalized eigenfunctions. Denote by \( \mathcal{D}'(\hat{\mathcal{C}}) \) the space of distributions on \( \hat{\mathcal{C}} \). We have a nuclear rigging (see [2], Section 14.2)
\[
\mathcal{D}(\hat{\mathcal{C}}) \subset L^2(\mathbb{C}, \mu_{a,b}) \subset \mathcal{D}'(\hat{\mathcal{C}}),
\]
and apply the usual formalism of generalized eigenfunctions, see [2], Chapter 15.

Recall that we have formally symmetric and formally commuting operators
\[
D_+ := \frac{1}{2}(\mathcal{D} + \overline{\mathcal{D}}), \quad D_- := \frac{1}{2i}(\mathcal{D} - \overline{\mathcal{D}})
\]
in \( L^2(\mathbb{C}, \mu) \) (defined on the domain \( \mathcal{D}(\hat{\mathcal{C}}) \)) and the spectral problem
\[
\mathcal{D} \Phi = \zeta \Phi, \quad \overline{\mathcal{D}} \Phi = \overline{\zeta} \Phi.
\]
Suppose that the operators $D_+, D_-$ admit commuting self-adjoint extensions. Then the operator $U$ of spectral decomposition can be written in terms of generalized eigenfunctions. Precisely, there exist a space $R$ equipped with a measure $\rho$ and an injective measurable map $r \mapsto \varphi_r$ from $R$ to $\mathcal{D}'(\hat{C})$ such that

$$D_+ \varphi_r = a(r) \varphi_r, \quad D_- \varphi_r = b(r) \varphi_r,$$

where $a(r), b(r)$ are real-valued functions, and the pairing

$$U f(r) = \{f, \varphi_r\}$$

of $f \in \mathcal{D}(\hat{C})$ and $\varphi_r$ determines a unitary operator $L^2(\mathbb{C}, \mu) \to L^2(R, \rho)$, see textbook [2], Subsect. 15.2.3.

Since the operator $\mathfrak{D}$ is elliptic, generalized eigenfunctions are smooth on $\hat{C}$, see e.g., [2], Theorem 16.2.1. Therefore in our case generalized eigenfunctions $\varphi_r$ are usual smooth solutions of the system of differential equations.

We also can identify the measure space $R$ with its image, and so we can think that the measure $\rho$ is sitting on the space $\Omega$ of smooth solutions of the systems (4.1), where $\zeta$ ranges in $\mathbb{C}$. We intend to show that for any measure $\rho$ on $\Omega$ the operator $J : L^2(\Omega, \rho) \to L^2(\mathbb{C}, \mu_{a,b})$ defined by

$$U h(z) = \int_{\Omega} h(r) \varphi_r(z) \, d\rho(r)$$

is not unitary. Precisely:

**Lemma 4.1.** Let $(a, b) \notin \Pi$. Let $\rho$ be a measure on $\Omega$, and let the corresponding operator $U$ be bounded. Then $\rho$ is an atomic measure supported by a finite set.

The idea of a proof is simple, it is explained in the next subsection, a formal proof is completed in Subsect. 4.3.

**Lemma 4.1** implies that for $(a, b) \notin \Pi$ the operators $D_+, D_-$ have no commuting self-adjoint extensions.

### 4.2. Almost proof of Lemma 4.1

For $\zeta$ being in a general position, the system (4.1) has a unique solution, and it has the form $2F_1^C[\cdot; z]$. Denote by $\Omega_{hyp}$ the subset of $\Omega$ consisting of the functions $2F_1^C[\cdot; z]$. We wish to prove the following statement:

**Lemma 4.2.** Let $(a, b) \notin \Pi$ and $a + b, a - b, a, b \notin \mathbb{Z}$. Let $\rho$ be a measure on $\Omega$, and let the corresponding operator $U$ be bounded. Then $\rho$ is atomic on $\Omega_{hyp}$.

**Proof.** Set $\zeta = \lambda^2$. Then a hypergeometric solution of the system (4.1) has one of the two forms:

$$2F_1^C \left[ \frac{a + \lambda(a - \overline{\lambda}, a - \lambda|a - \overline{\lambda})}{a + b|a + b}; z \right], \quad 2F_1^C \left[ \frac{a + \lambda(a + \overline{\lambda}, a - \lambda|a + \overline{\lambda})}{a + b|a + b}; z \right].$$

In the first case we have $(a + \lambda) - (a - \overline{\lambda}) = 2 \Re \lambda \in \mathbb{Z}$, hence $\lambda = \frac{1}{2}(k + is)$, where $k \in \mathbb{Z}$, $s \in \mathbb{R}$. We come to the functions $\mathcal{K}_{a,b}(z; k, is)$.

In the second case we have $\lambda - \overline{\lambda} \in \mathbb{Z}$, i.e., $\lambda = \tau \in \mathbb{R}$. We come to the functions (4.2)

$$\mathcal{K}(z; 0, \tau) = 2F_1^C \left[ \frac{a + \tau(a + \tau, a - \tau|a - \tau)}{a + b|a + b}; z \right].$$

---

[21] Basically, this is a result of Kostyuchenko and Mityagin [23]–[24] with weaker conditions for a rigging.
Next, we will show that

\[(4.3) \quad \text{the measure } \rho \text{ is zero on the set of all } \lambda = \frac{1}{2}(k + is) \text{ with } s \neq 0.\]

Our kernel has the following asymptotics at \(z = 0\) and \(z = 1:\)

\[(4.4) \quad \mathcal{K}(z; k, is) = (1 + O(z)) + B(k, is)|z|^{2-2a-2b}(1 + O(z)) \quad \text{as } z \to 0,
\]

\[\mathcal{K}(z; k, is) = C(k, is)(1 + O(1 - z)) + D(k, is)|1 - z|^{2b-2a}(1 + O(1 - z)) \quad \text{as } z \to 1,\]

where the coefficients \(B, C, D\) are continuous non-vanishing functions on \(\Lambda\) and all \(O(\cdot)\) are uniform on compact subsets of \(\Lambda\) (see formulas (3.20)–(3.31)).

For definiteness, assume that \(a + b > 2\). Consider a point \((k_0, is_0) \in \Lambda, s_0 \neq 0\) and a neighborhood \(N\) of \((k_0, is_0)\). Assume that \(\rho(N) > 0\). Denote by \(I_N\) the indicator function of the set \(N\). The function \(U_I\) has the following asymptotics at \(z = 0:\)

\[UI_N(z) = \alpha (1 + O(z)) + \beta |z|^{2-2a-2b}(1 + O(z)) \quad \text{as } z \to 0.\]

Due to uniformity \(O(\cdot)\), for a sufficiently small neighborhood \(N\) we have \(\alpha \neq 0, \beta \neq 0\). Since \(a + b > 2\), the actual asymptotics is

\[UI_N(z) = \beta |z|^{2-2a-2b}(1 + O(z)).\]

Therefore

\[UI_N \notin L^2(\mathbb{C}, |z|^{2a+2b-2}|1 - z|^{2a-2b} d\mathfrak{v}).\]

This contradicts to boundedness of \(U\). Thus any point has a neighborhood of zero measure, and this implies claim (4.3) in the case \(a + b > 1\).

In domains \(a + b < 0, a - b < 1, a - b > 1\) we get the same effect.

Next, examine the complementary series \(\mathcal{K}(z; 0, \tau)\) of eigenfunctions, see (4.2). We have the same asymptotics (4.4)–(4.5), we only must write the coefficients of the form \(A(0, \tau), B(0, \tau), C(0, \tau), D(0, \tau)\) in (4.4)–(4.5). These functions have zeros and poles on the axis \(\tau \in \mathbb{R}\). The same argument as above shows that if \(\tau_0\) is not a zero and not a pole of all our coefficients, then the measure \(\rho\) is zero on a sufficiently small neighborhood of \(\tau_0\). The set of zeros and poles is countable. This completes the proof of the lemma. \(\square\)

4.3. Proof of non-self-adjointness. However, our system of differential equations (4.1) has solutions that have not the form \(2F_1^C\), and enumeration of all possible degenerations is tedious. So we continue the proof of Lemma 4.1 without constrains of Lemma 4.2. Due to the homographic transformations, without loss of generality we can set

\[(4.6) \quad a + b > 2.\]

First, we examine asymptotics in a neighborhood of \(z = 0\).

Asymptotics at \(z = 0\). Non-logarithmic case. If \(a + b \neq 2, 3, \ldots\), then the equation \(\mathcal{D}\Phi = \lambda^2 \Phi\) has two holomorphic solutions,

\[\Psi_1(z) := 2F_1\left[\begin{array}{c} a + \lambda, a - \lambda, a + b \end{array}; 1 - b + \lambda, 1 - b - \lambda, 2 - a - b \end{array}; z\right], \quad \Psi_2(z) := z^{1-a-b} 2F_1\left[\begin{array}{c} 1 - b + \lambda, 1 - b - \lambda, 2 - a - b \end{array}; 2 - a - b \end{array}; z\right].\]

The equation \(\overline{\mathcal{D}} \Phi = \overline{\lambda}^2 \Phi\) has two antiholomorphic solutions

\[\Psi_1(\overline{\tau}) := 2F_1\left[\begin{array}{c} a + \overline{\lambda}, a - \overline{\lambda}, a + b \end{array}; \overline{\tau}\right], \quad \Psi_2(\overline{\tau}) := \overline{\tau}^{1-a-b} 2F_1\left[\begin{array}{c} 1 - b + \overline{\lambda}, 1 - b - \overline{\lambda}, 2 - a - b \end{array}; 2 - a - b \end{array}; \overline{\tau}\right].\]
Therefore a single-valued solution of the system must have the form
\[
A\Psi_1(z) \overline{\Psi}_1(z) + B\Psi_2(z) \overline{\Psi}_2(z).
\]
The first term has \(L^2(\mathbb{C}, \mu_{a,b})\)-asymptotics at \(z = 0\), by (4.6) the second term
has non-\(L^2\)-asymptotics. Thus the spectral measure \(\rho\) is supported by the set of
functions of the form \(\Psi_1(z) \overline{\Psi}_1(z)\).

**Asymptotics at** \(z = 0\). **Logarithmic case.** Now let \(a + b = n = 2, 3, \ldots\).
Then the equation \(\mathcal{D} \Phi = \lambda^2 \Phi\) has two holomorphic solutions,

\[
\Psi_1(z) = 2F_1[a + \lambda, a - \lambda; n, z], \quad \Psi_2(z),
\]
where \(\Psi_2(z)\) is a logarithmic solution, which has the form (3.39). The equation \(\mathcal{D} \Phi = \overline{\lambda}^2 \Phi\) has two antiholomorphic solutions,

\[
\overline{\Psi}_1(z) = 2F_1[a + \overline{\lambda}, a - \overline{\lambda}; n, \overline{z}], \quad \overline{\Psi}_2(z).
\]

A single valued solution must have the form
\[
A\Psi_1(z) \overline{\Psi}_1(z) + B\left(\Psi_1(z) \overline{\Psi}_2(z) + \Psi_2(z) \overline{\Psi}_1(z)\right).
\]
The asymptotics of the second summand is \((z^{-n+1} + \overline{z}^{-n+1}) + O(z^{-n+2})\) if \(n \geq 3\).
If \(n = 2\) we have \((z^{-1} + \overline{z}^{-1}) + O(z^{-\varepsilon})\). We get a non-\(L^2\) asymptotics.

Thus, for \(a + b > 2\) the spectral measure is supported by set of functions of the
form \(\Psi_1(z) \overline{\Psi}_1(z)\).

**Single-valuedness near** \(z = 1\). **Non-logarithmic case.** Assume that \(a - b \notin \mathbb{Z}\). We apply formulas Erdélyi, et al., [9], (2.10.1), (2.10.5) and write explicit expansions of \(\Psi_1, \overline{\Psi}_1\) at \(z = 1\).

\[
\Psi_1(z) = A_1 G_1(1 - z) + A_2(1 - z)^{b-a} G_2(1 - z);
\overline{\Psi}_1(z) = A_1 \overline{G}_1(1 - \overline{z}) + A_2(1 - \overline{z})^{b-a} \overline{G}_2(1 - \overline{z}),
\]
where \(G_1, G_2\) are certain series \(2F_1\) and the coefficients \(A_1, A_2\) are products
of gamma functions, see the explicit formulas \((3.42)-(3.43)\) above. Clearly, the product
\(\Psi_1(z) \overline{\Psi}_1(z)\) can be single-valued only if \(A_2 = \overline{A}_2 = 0\), or \(A_1 = \overline{A}_1 = 0\). Looking
to the explicit expressions for the gamma-coefficients, we observe that the first case
happens if both hypergeometric series \(G_1(z), G_2(z)\) are terminating (i.e., \(a - \lambda = 0, -1, \ldots\) or \(a + \lambda = 0, -1, \ldots\), in particular, \(\lambda\) is real). The second variant holds if only if both series \(G_2(1 - z), \overline{G}_2(1 - \overline{z})\) are terminating (i.e., \(b - \lambda = 0, -1, \ldots\) or \(b + \lambda = 0, -1, \ldots\)).

**Single-valuedness near** \(z = 1\). **Logarithmic case.** Now let \(b - a \in \mathbb{Z}\). The transposition \(a \leftrightarrow b\) corresponds to a homographic transformation of differential
operators, it preserves the condition \(a + b \geq 2\). Therefore we can assume \(m := b - a \geq 0\). Represent \(\Psi_1(z), \overline{\Psi}_1(z)\) as combinations of basic solutions of the hypergeometric
equations at the point \(z = 1\),

\[
\Psi_1(z) = A_2 F_1[a + \lambda, a - \lambda; b - a + 1; z] + B \Theta(1 - z);
\overline{\Psi}_1(z) = \overline{A}_2 F_1[a + \overline{\lambda}, a - \overline{\lambda}; b - a + 1; \overline{z}] + \overline{B} \overline{\Theta}(1 - \overline{z}),
\]
where \(\Theta(1 - z)\) is a logarithmic series of the type (3.39), see Erdélyi, et al., [9],
(2.10.12). A straightforward calculation shows that the product \(\Psi_1(z) \overline{\Psi}_1(z)\) can
be single valued near \(z = 1\) only if \(B = \overline{B} = 0\). Therefore \(\Psi_1(z)\) is single valued
near \(z = 1\), and therefore it is a single valued solution of a hypergeometric equation.
on the whole plane $\hat{\mathbb{C}}$. Hence (see Erdélyi, et al., [9], Subsect. 2.2.1) $\Psi_1(z)$ is a polynomial.

**Behavior at infinity.** Thus the spectral measure $\rho$ is supported by generalized eigenfunctions of the following types

$$p_1(z)p_2(\bar{z}), \quad (1 - z)^{b-a|b-a} q_1(z)q_2(\bar{z}),$$

where $p_j$, $q_j$ are polynomials. However, our density $\mu_{a,b}(z)$ has a behavior $\sim |z|^{4a-2}$ at infinity and therefore the space $L^2$ can contain only a finite number orthogonal functions of such a type. \qed

5. **Symmetry of differential operators**

Here we show that $J_{a,b}^*$ sends $\mathcal{D}_{\text{even}}(\hat{\Lambda})$ to $\mathcal{R}_{a,b}$ and verify that $\mathcal{D}$ and $\overline{\mathcal{D}}$ are adjoint one to another on $\mathcal{R}_{a,b}$.

In this section we denote by $D_r(u) \subset \mathbb{C}$ (resp. $\overline{D}_r(u)$) the open (resp. closed) disc in $\mathbb{C}$ of radius $r$ with center at $u$. By $S_r(u)$ we denote the circle $|z - u| = r$.

5.1. **The map $J_{a,b}^*$ on the space $\mathcal{D}_{\text{even}}(\hat{\Lambda})$.** Introduce a natural topology in the space $\mathcal{R}_{a,b}(\hat{\mathbb{C}})$ defined in Subsect. 1.1. Consider the space $\mathcal{R}(0)$ of functions in $\overline{D}_{1/3}(0)$ having the form $\alpha(z) + \beta(z)|z|^{2a+2b-2}$, where $\alpha(z)$, $\beta(z)$ are smooth in $\overline{D}_{1/3}(0)$ up to the boundary. Let $C_{\text{flat}}^\infty(\overline{D}_{1/3}(0)) \subset C^\infty(\overline{D}_{1/3}(0))$ be the subspace consisting of all functions that are flat at 0. The space $\mathcal{R}(0)$ is a quotient space

$$\mathcal{R}(0) \simeq \left[ C^\infty(\overline{D}_{1/3}(0)) \oplus |z|^{2a+2b-2}C^\infty(\overline{D}_{1/3}(0)) \right] / C_{\text{flat}}^\infty(\overline{D}_{1/3}(0)).$$

We equip $\mathcal{R}(0)$ with the topology of a quotient space. In the same way we define a topology in the space $\mathcal{R}(1)$ of smooth functions in $\overline{D}_{1/3}(1)$ having the form $\gamma(z) + \delta(z)|1 - z|^{2a-2b}$.

We define a topology in $\mathcal{R}_{a,b}$ as a weakest topology such that:

- a) The restriction operators

$$\mathcal{R}_{a,b} \rightarrow \mathcal{R}(0), \quad \mathcal{R}_{a,b} \rightarrow \mathcal{R}(1), \quad \mathcal{R}_{a,b} \rightarrow C^\infty(\overline{D}_2(0) \setminus (D_{1/3}(0) \cap D_{1/3}(1)))$$

are continuous.

- b) For all $\alpha$, $\beta$, $N$ the following seminorms are continuous

$$(5.1) \quad p_{\alpha,\beta,N}(f) = \sup_{C \setminus D_2(0)} |z|^{2+\alpha+\beta} |\ln|z||^N \left| \frac{\partial^{\alpha+\beta} f(z)}{\partial z^\alpha \partial \overline{z}^\beta} \right|.$$

Recall that $\hat{\Lambda} := \Lambda \setminus \{(0,0)\}$.

**Lemma 5.1.** For $|z| > 2$, $(k,s) \in \hat{\Lambda}$ we have the following expansion

$$(5.2) \quad \mathcal{K}(z; k, is) = z^{-\alpha - \frac{k+i}{s}}e^{-\frac{k+i}{s}}B(k, s; z^{-1}) + z^{-\alpha - \frac{k+i}{s}}e^{\frac{k+i}{s}}B(-k, -s; z^{-1}),$$

where the expression $B(k, s; u)$ for fixed $k$ is smooth $s$ except the point $(k, s) = (0,0)$.

**Proof.** We refer to expansion (3.32)–(3.33). Notice that for $k = 0$, $s = 0$ we have a singularity in this expansion (but the kernel itself is analytic at this point). \qed
Proposition 5.2. a) Let $\Phi \in \mathcal{D}_{\text{even}}(\hat{\Lambda})$. Then $J^*_{a,b} \Phi \in \mathcal{R}_{a,b}$.

b) Moreover, the operator $J^*_{a,b}$ is a continuous operator from $\mathcal{D}_{\text{even}}(\hat{\Lambda})$ to $\mathcal{R}_{a,b}$.

**Proof.** Forms of asymptotics of $J^*_{a,b} \Phi$ at 0 and 1 follow from the expressions \[ (3.20), (3.29) \]. Let us examine the asymptotics at $z \to \infty$. Without loss of generality we can assume that $|k|$ is fixed. We write

$$J^*_{a,b} \Phi(z) = z^{-a-b} |a+b|^{-\frac{1}{2}} \int_{\mathbb{R}} z^{-\frac{1}{2}} B(k, s; z^{-1}) \Phi(k, s) \, ds + \{ \text{similar term} \}$$

Differentiating the first summand by $\frac{\partial^{\alpha+\beta}}{\partial z^{\alpha} \partial z^*^{\beta}}$ and keeping in mind \[ (3.32) \] and Lemma 3.6, we get an expression of the form

$$z^{-a-b-\frac{1}{2}-a+b+\frac{1}{2}} \sum_{0 \leq p \leq \alpha, 0 \leq q \leq \beta} \int_{\mathbb{R}} z^{-\frac{1}{2}} U_{p,q}^{a,\beta}(a, b, k, s) \times \frac{\Gamma^C(-k-is|k+is)(a+k+is)_p(a-k-is)_q(a+k-is)_q(a-k-is)_q \times \Gamma^C(b-\frac{k+is}{2}|b-\frac{k-is}{2}) \Gamma^C(a+k-is|a+k-is)_p(a+b)_q \times 2F_1 \left[ a+\frac{k+is}{2} + p, a+\frac{k-is}{2} + p; z^{-1} \right] \times 2F_1 \left[ a+\frac{k-is}{2} + q, a+\frac{k-1}{2} + q; z^{-1} \right] \times \Phi(k, s) \, ds,$$

where $U_{p,q}^{a,\beta}(a, b, k, s)$ are polynomials. It is easy to verify that the integrand is a smooth compactly supported function on $\hat{\Lambda}$. Next, we write

$$|z|^{-is} = i \frac{\partial}{\partial s} \ln |z|^{-is},$$

integrate our expansion by parts $N$ times and observe that $p_{a,\beta,N}(J^*_{a,b} \Phi) \to \infty$.

The continuity follows from the same considerations. \hfill \Box

As a corollary, we obtain the following lemma.

**Lemma 5.3.** The operator $J^*_{a,b}$ is continuous as an operator from $\mathcal{D}_{\text{even}}(\hat{\Lambda})$ to the space $L^2(\mathbb{C}, \mu_{a,b})$.

**Proof.** Indeed, for $(a, b) \in \Pi$ the identical embedding $f \mapsto f$ of $\mathcal{R}_{a,b}$ to $L^2(\mathbb{C}, \mu_{a,b})$ is continuous. \hfill \Box

**Lemma 5.4.** If $f \in \mathcal{R}_{a,b}$, then $\mathcal{D} f \in \mathcal{R}_{a,b}$.

**Proof.** Let us check the behavior of $\mathcal{D} f$ at 0. For definiteness assume that $a+b \neq 1$. Then near zero we have

$$\mathcal{D} f = \mathcal{D} (\alpha(z) + \beta(z)|z|^{1-a-b}) = \left\{ \left( z(1-z) \frac{d^2}{dz^2} + (a+b) \frac{d}{dz} \right) z^{1-a-b} \right\} \cdot \alpha(z) \cdot \beta(z) + \{ \text{the rest} \}.$$ 

Obviously, the rest has the form $\alpha(z) + \beta(z)|z|^{2-a-b}$ with smooth $\alpha, \beta$. The expression in the curly bracket is $-(a+b)(a+b-1)z^{1-a-b}$. \hfill \Box

5.2. symmetry of differential operators.

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22Cf. [11], Sect.1.2.
Proposition 5.5. For any $f, g \in \mathcal{R}_{a,b}(\mathbb{C})$
\[ \langle \mathcal{D} f, g \rangle = \langle f, \overline{g} \rangle. \]

Proof. Let $f, g \in \mathcal{R}_{a,b}$. We wish to show that
\[ \int_{\mathbb{C}} (\mathcal{D} f(z) \cdot g(z) - f(z) \cdot \overline{g(z)}) \mu_{a,b} \, dz = 0. \]

By Lemma 5.3, $\mathcal{D} f, \mathcal{D} g \in L^2(\mathbb{C}, \mu_{a,b})$. Therefore our improper integral absolutely converges, we write it as
\[ \lim_{\varepsilon \to 0} \int_{D_{1/\varepsilon}(0) \setminus (D_{\varepsilon}(0) \cup D_{\varepsilon}(1))} \cdots \frac{dz \wedge d\overline{z}}{2i}. \]

Next, we integrate two times by parts in $z$ (with the Green formula) and after a simple calculation come to
\[ \lim_{\varepsilon \to 0} \left\{ \int_{S_{1/\varepsilon}(0)} V(z) d\overline{z} - \int_{S_{\varepsilon}(0)} V(z) d\overline{z} - \int_{S_{\varepsilon}(1)} V(z) d\overline{z} \right\}, \]

where
\[ V(z) = \left( \frac{\partial f}{\partial z} \cdot g(z) - f(z) \cdot \frac{\partial g(z)}{\partial z} \right) \overline{z}(1 - z) \mu_{a,b}(z). \]

We claim that all summands in (5.3) tend to 0. For the first summand this is clear. For the second summand we represent $f, \overline{g}$ as
\[ f(z) = \alpha(z) + \beta(z) z^{1-a-b} 1-a-b, \quad \overline{g}(z) = \gamma(z) + \delta(z) z^{1-a-b} 1-a-b. \]

Then $V(z)$ transforms to an expression of the following type:
\[ \left( A(z) + B(z) z^{-a-b} 1-a-b + C(z) z^{1-2a-2b} -2-2a-2b \right) \times \]
\[ \times z(1-z) \cdot z^{a+b-1} a+b-1 (1-z)^{a-b} \]
where $A(z), B(z), C(z)$ are smooth near 0. We emphasize that the term with $z^{1-2a-2b} -2-2a-2b$ in the bracket appears with the coefficient
\[ (2-2a-2b) (\beta(z) \delta(z) - \beta(z) \delta(z)) = 0. \]

Thus we get summands with the following behavior at 0:
\[ \sim A(0) z^{a+b} a+b-1, \quad \sim B(0) z^0, \quad \sim C(0) z^{2-a-b} 1-a-b. \]

Since $0 < a + b < 2$ all powers are $>-1$ and therefore\[ \int_{|z|=\varepsilon} \cdots d\overline{z} \]tends to 0.

□

6. The operator $J_{a,b}^*$ is an isometry

Here we prove half of Theorem 1.3

6.1. The statement. First, denote by $\Lambda_+$ the subset of $\Lambda$ consisting of $(k+i\varepsilon)/2$ such that $k > 0$ or $k = 0$ and $s > 0$. We have an obvious identification $\mathcal{D}_{\text{even}}(\Lambda) \simeq \mathcal{D}(\Lambda_+)$. Let $u(\lambda), v(\lambda)$ be smooth compactly supported function on $\Lambda_+$. Then
\[ \langle J_{a,b}^* u, J_{a,b}^* v \rangle_{L^2(\mathbb{C}, \mu_{a,b})} = 2 \langle u, v \rangle_{L^2(\Lambda_+, \kappa_{a,b})}. \]

\[ \text{See a discussion of a parallel situation for ordinary differential operators in [36], Section 1. However, in the one-dimensional case we must impose boundary conditions in such points.} \]
Our proof is based on heuristic arguments outlined in Berezin, Shubin \[3\], Section 2.6, for ordinary differential operators. However, this way is tiresome.

### 6.2. Preliminary remarks.

Recall that
\[ J_{a,b}^* u(z) = 2 \int_{\Lambda^+} u(\lambda) \mathcal{K}(z,\lambda) \, \zeta_{a,b}(\lambda) d\lambda. \]

By Lemma 5.3 this operator is continuous as an operator $\mathcal{D}(\Lambda^+) \to L^2(\hat{\mathbb{C}},\mu_{a,b})$. Therefore the sesquilinear form
\[ (6.1) \]
\[ T(u, v) := \langle J_{a,b}^* u, J_{a,b}^* v \rangle_{L^2(\hat{\mathbb{C}},\mu_{a,b})} \]
is continuous as a form $\mathcal{D}(\Lambda^+) \times \mathcal{D}(\Lambda^+) \to \mathbb{C}$. By the kernel theorem (see, e.g., \[19\], Sect. 5.2) it is determined by a distribution. Formally, we transform (6.1) as
\[ (6.2) \]
\[ \int_{\mathbb{C}} \int_{\Lambda^+} u(\lambda) \mathcal{K}(z,\lambda) \zeta_{a,b}(\lambda) \, d\lambda d\bar{\lambda} = \int_{\Lambda^+} \int_{\Lambda^+} u(\lambda) \mathcal{K}(z,\lambda) \zeta_{a,b}(\lambda) \zeta_{a,b}(\nu) \, d\lambda d\nu, \]

where
\[ (6.3) \]
\[ H(\lambda,\nu) = \int_{\mathbb{C}} \mathcal{K}(z,\lambda) \mathcal{K}(z,\nu) \mu_{a,b}(z) \, d\mathbb{B}. \]

Notice that all integrals in line (6.2) converge absolutely. However, the triple integral $\int_{\Lambda^+} \int_{\Lambda^+} \int_{\mathbb{C}}$ is not absolutely convergent. The integrand in (6.4) decreases as $|z|^{-2}$ and the integral diverges.

However, we regard $H(\lambda,\nu)$ as a distribution, then Lemma 6.1 can be reformulated in the form:

**Lemma 6.2.** We have the following identity of distributions on $\mathcal{D}(\Lambda^+) \times \mathcal{D}(\Lambda^+)$:
\[ (6.5) \]
\[ H(\lambda,\nu) = \delta(\lambda - \nu). \]

### 6.3. Orthogonality of packets.

**Lemma 6.3.** Let $u, v \in \mathcal{D}(\Lambda^+)$ and supports $\text{supp}(u), \text{supp}(v)$ have empty intersection. Then
\[ \langle J_{a,b}^* u, J_{a,b}^* v \rangle_{L^2(\hat{\mathbb{C}},\mu_{a,b})} = 0. \]

**Proof.** Denote $D_+ := \frac{1}{2}(\mathbb{D} + \overline{\mathbb{D}})$, $D_- := \frac{1}{2}(\mathbb{D} - \overline{\mathbb{D}})$. By Proposition 5.2 $J_{a,b}^* u$ is contained in the space $\mathcal{R}_{a,b}$. By Proposition 6.3 the operators $D_+, D_-$ are formally symmetric on $\mathcal{R}_{a,b}$. Since they formally commute, for any real polynomial $p(D_+, D_-)$ we have
\[ (p(D_+, D_-) J_{a,b}^* u, J_{a,b}^* v) = (J_{a,b}^* u, p(D_+, D_-) J_{a,b}^* v), \]
or
\[ (6.6) \]
\[ \langle J_{a,b}^* p(\text{Re} \, \lambda, \text{Im} \, \lambda) \cdot u, J_{a,b}^* v \rangle = \langle J_{a,b}^* u, J_{a,b}^* p(\text{Re} \, \lambda, \text{Im} \, \lambda) \cdot v \rangle, \]

where $\cdot$ denotes the operator of multiplication by a function. We choose a sequence $p_N$ of polynomials such that $p_N$ uniformly converges to 1 on $\text{supp}(u)$ with all derivatives and converges to 0 on $\text{supp}(v)$. By Lemma 5.3 the map $J_{a,b}^*$ is continuous as a map $\mathcal{D}(\Lambda^+) \to L^2(\mathbb{C},\mu_{a,b})$. Replacing $p$ by $p_N$ in (6.6) and passing to a limit, we come to the desired statement. \[\square\]
6.4. Next reduction of our statement. Let $S(u, v)$ be an Hermitian form on $\mathcal{D}(\Lambda_+)$. We say that $S$ is $C^\omega$-smooth if it has the form

$$S(u, v) = \int_{\Lambda_+} \int_{\Lambda_+} M(\lambda, \nu) u(\lambda) \overline{v(\nu)} \, d\lambda \, d\nu,$$

where $M$ is a real analytic function on $\Lambda_+ \times \Lambda_+$.

Lemma 6.4. We have

$$\langle J^*_a b u, J^*_a b v \rangle_{L^2(\mathbb{C}, \mu_{a,b})} = \langle u, v \rangle_{L^2(\Lambda, \kappa_{a,b})} + S(u, v),$$

where $S(u, v)$ is $C^\omega$-smooth.

This lemma together with Lemma 6.3 imply the desired statement, i.e., the identity (6.5). Indeed, for any $u, v$ with disjoint support, we have

$$\int_{\Lambda_+} \int_{\Lambda_+} M(\lambda, \nu) u(\lambda) \overline{v(\nu)} \kappa_{a,b}(\lambda) \kappa_{a,b}(\nu) \, d\lambda \, d\nu = 0,$$

therefore $M(\lambda, \nu) = 0$.

The rest of this section is occupied by the proof of Lemma 6.4.

6.5. Beginning of the proof of Lemma 6.4 Cleaning of the problem. Step 1. Represent

$$u = \sum_k u_k \delta(\text{Re} \lambda - k/2), \quad v = \sum_l v_l \delta(\text{Re} \lambda - l/2),$$

in fact the sums are finite and $u_k, v_l$ depend on a real variable $s$. By Lemma 6.3 we have

$$\langle J^*_a b u_k, J^*_a b v_l \rangle = 0 \quad \text{for } k \neq l.$$ 

Therefore it is sufficient to examine only inner products

$$\langle J^*_a b u_k, J^*_a b v_k \rangle = \int \mathcal{R}(z) \, d\overline{z},$$

where

$$\mathcal{R}(z) := \int_{\Lambda_+} u_k(\text{is}) \mathcal{K}(z, \frac{1}{2}(k + is)) \kappa_{a,b}(\frac{1}{2}(k + is)) \, ds \times \int_{\Lambda_+} \overline{v_k(it)} \mathcal{K}(z, \frac{1}{2}(k + it)) \kappa_{a,b}(\frac{1}{2}(k + it)) \, dt \, \mu_{a,b}(z).$$

Step 2. Represent the integral as $\int_{|z| < 2} \mathcal{R} + \int_{|z| > 2} \mathcal{R}$.

Let us show that the first summand is $C^\omega$-smooth. In this case the triple integral absolutely converges and can be written as

$$\int_{|z| < 2} \mathcal{R} \, d\overline{z} = \int_{\mathbb{R}} \int_{\mathbb{R}} u_k(\text{is}) \overline{v_k(it)} L(s, t) \, ds \, dt,$$

where

$$L(s, t) = \int_{|z| < 2} \mathcal{K}(z, \frac{1}{2}(k + is)) \mathcal{K}(z, \frac{1}{2}(k + it)) \mu_{a,b}(z) \, d\overline{z}.$$ 

Integrand makes sense for complex $s, t$ that are sufficiently close to $\mathbb{R}$ and the integral absolutely converges (singularities at $z = 0$ and 1 have the forms (4.4), (4.5)). Therefore $L(s, t)$ is a holomorphic function in $s, t$ near $\mathbb{R} \times \mathbb{R}$. 

Therefore our question is reduced to an examination the integral
\[ \int_{|z| > 2} R(z) \, d\overline{z} \]

Step 3. A decomposition of the kernel. Applying Theorem 3.9.c, we write \( \mathcal{K}(z, \lambda) \) in the domain \( |z| > 2 \) as

\[ \mathcal{K}(z, \lambda) = W_1 + W_2 + W_3 := \]
\[ = A(\lambda)(-z)^{a-\lambda-\overline{\lambda}} + A(-\lambda)(-z)^{-a+\lambda-\overline{\lambda}} + \Psi(z, \lambda), \]

where
\[ A(\lambda) = \frac{\Gamma^C(2\lambda - 2\overline{\lambda})}{\Gamma^C(b - \lambda)b + \lambda\Gamma^C(a - \lambda|a + \overline{\lambda})} \]

and
\[ \Psi(z, \lambda) = O(|z|^{-2a-1}) \quad \text{as} \quad z \to \infty. \]

Notice that
\[ |A(\lambda)|^2 = A(\lambda) A(-\lambda) = \kappa_{a,b}^{-1}(\lambda). \]

Therefore the integral \( \int_{|z| > 2} R(z) \, d\overline{z} \) splits into a sum of 9 summands \( V_{\alpha\beta} \), where \( \alpha, \beta = 1, 2, 3 \),

\[ V_{\alpha\beta} := \int_{|z| > 2} \int_{\mathbb{R}} W_\alpha(z; k, s) u_k(\mu) \kappa_{a,b}(\frac{1}{2}k + \overline{\mu}) \, ds \times \]
\[ \quad \times \int_{\mathbb{R}} W_\beta(z; k, t) v_k(it) \kappa_{a,b}(\frac{1}{2}k + it) \, dt \cdot \mu_{a,b}(z) \, d\overline{z}. \]

Step 4. For five summands \( V_{13}, V_{23}, V_{31}, V_{32}, V_{33} \) we immediately get absolute convergence of triple integrals and \( C^\omega \)-smoothness. For instance,

\[ V_{13} = \int_{\mathbb{R}} \int_{\mathbb{R}} u_k(\mu) v_k(it) A(\frac{1}{2}k - \mu) \left( \frac{\mu - z}{\overline{\mu} - z} \right)^{-k} \times \]
\[ \quad \times \left[ \int_{|z| > 2} \left( \frac{\mu}{\overline{\mu}} \right)^{-k} |z|^{-2a+is} \Psi(z, \frac{1}{2}k + is) \mu_{a,b}(z) \, d\overline{z} \right] \, ds \, dt. \]

We simplified the integrand using (6.9)). The expression in the square brackets is real analytic (the integrand decreases as \( |z|^{-3} \)).

Step 5. Non-obvious summands are \( V_{11}, V_{12}, V_{21}, V_{22} \). We start with \( V_{11} \),

\[ V_{11} := \int_{|z| > 2} \int_{\mathbb{R}} u_k(\mu) A(\frac{1}{2}k + \overline{\mu}) \left( \frac{\mu}{z} \right)^{-k/2} |z|^{-2a+is} \kappa_{a,b}(\frac{1}{2}k + is) \, ds \times \]
\[ \quad \times \int_{\mathbb{R}} A(\frac{1}{2}k + it) v_k(it) \left( \frac{z}{\overline{\mu}} \right)^{k/2} |z|^{-2a-it} \kappa_{a,b}(\frac{1}{2}k + it) \, dt \mu_{a,b}(z) \, d\overline{z}. \]

For \( k = 0 \) we must keep in mind that the integration \( \int_{\mathbb{R}} \) actually is taken over a ray \([\varepsilon, \infty)\) for some \( \varepsilon > 0 \). Applying (6.9), we come to

\[ V_{11} := \int_{|z| > 2} \int_{\mathbb{R}} u_k(\mu) v_k(it) A(\frac{1}{2}k - \mu) A(\frac{1}{2}k + it) \left( \frac{\mu}{z} \right)^{-1} |z|^{-4a+is+it} \, ds \, dt \times \]
\[ \quad \times \mu_{a,b}(z) \, d\overline{z}. \]
Next, we notice that
\[ \mu_{a,b}(z) = |z|^{2a+2b-2} |1 - z|^{2a-2b} = |z|^{4a-2} + O(|z|^{4a-3}) \quad \text{as } z \to \infty. \]

We write
\[ (6.11) \quad \mu_{a,b}(z) = |z|^{4a-2} + (\mu_{a,b}(z) - |z|^{4a-2}), \]

substitute this to (6.10) and decompose (6.10) as a sum of two integrals. The second summand immediately gives a \( C^\omega \)-smooth term. The first summand is the topic of our interest. It equals the following expression:
\[ (6.12) \quad I(u, v) := \int_{|z|>2} \int_{\mathbb{R}} \int_{\mathbb{R}} u_k(is)v_k(it)A(\frac{1}{2}(k - is))^{-1} A(\frac{1}{2}(k + it))^{-1} \times \]
\[ \times |z|^{-2-is+it} ds \, dt \, d\mathbb{F}. \]

### 6.6. Application of the Sokhotski formula and disappearance of a singular term.

**Step 6. Extension to the complex domain.** Now consider a function \( I(u, v, \varepsilon) \) obtained by replacing \( s \to s - i\varepsilon \) in the boxed term, \( \varepsilon > 0 \). The new triple integral absolutely converges, we can change the order of integrations and explicitly integrate in \( z \). We get
\[ I(u, v, \varepsilon) = \int_{\mathbb{R}} \int_{|z|>2} u_k(is)v_k(it)A(\frac{1}{2}(k - is))^{-1} A(\frac{1}{2}(k + it))^{-1} \frac{2^{-is-\varepsilon+it}}{-is - \varepsilon + it} ds \, dt. \]

Next, we claim that
\[ I(u, v) = \lim_{\varepsilon \to +0} I(u, v, \varepsilon). \]

Indeed, we integrate \( I(u, v, \varepsilon) \) two times by parts in \( s \) and come to
\[ I(u, v, \varepsilon) = \int_{|z|>2} \int_{\mathbb{R}} \frac{\partial^2}{\partial s^2} \left[ u_k(is)A(\frac{1}{2}(k - is))^{-1} \right] \times \]
\[ \times |z|^{-2-is-\varepsilon+it} ds \, dt \, d\mathbb{F}. \]

The new triple integral absolutely converges and is continuous at \( \varepsilon = +0 \).

Thus we come to the so-called distribution \( \frac{1}{x - i\varepsilon} \), see, e.g., [13]. Recall the Sokhotski formula
\[ (6.13) \quad \lim_{\varepsilon \to +0} \int_{\alpha}^{\beta} \frac{f(y) \, dy}{x - y - i\varepsilon} = \text{p.v.} \int_{\alpha}^{\beta} \frac{f(y) \, dy}{x - y} + \pi i f(x), \]

where p.v. denotes the principal value of an integral.

Applying this formula and keeping in mind (6.9), we come to
\[ (6.14) \quad I(u, v) = \text{p.v.} \int \int u_k(is)v_k(it)A(\frac{1}{2}(k - is))^{-1} A(\frac{1}{2}(k + it))^{-1} \frac{2^{-is+it}}{-is + it} ds \, dt + \]
\[ + \pi \int u_k(is)v_k(is) \varepsilon\mu_{a,b}(\frac{1}{2}(k + is)) ds. \]
Step 8. We deal with $V_{22}$ in the same way and come to

\begin{equation}
V_{22} = \text{p.v.} \int \int_{\mathbb{R}} u_k(is) v_k(it) A\left(\frac{1}{2}(k - is)\right)^{-1} A\left(\frac{1}{2}(k + it)\right)^{-1} \frac{2^{-it+is}}{-it + is} ds \, dt + \pi \int_{\mathbb{R}} u_k(is) v_k(is) \varphi_{a, b} \left(\frac{1}{2}(k + is)\right) ds + \left\{ \text{a } C^\omega\text{-smooth term} \right\}.
\end{equation}

Next, we take the sum $V_{11} + V_{22}$ modulo $C^\omega$-smooth terms. The expression

\begin{align}
A\left(\frac{1}{2}(k - is)\right)^{-1} A\left(\frac{1}{2}(k + it)\right)^{-1} & 2^{-it+is} - A\left(\frac{1}{2}(k - it)\right)^{-1} A\left(\frac{1}{2}(k + is)\right)^{-1} 2^{-it+is} \\
& - i(s - t)
\end{align}

has the form

\[ L(t, s) - L(s, t) \]

with analytic $L(t, s)$. It has a removable singularity on the line $t = s$. Thus the first summands in (6.14) and (6.15) give us a $C^\omega$-smooth term, the second summands give us the first term in (6.7), i.e., the desired delta-function.

6.7. End of the proof of Lemma 6.4

Step 9. Next, we examine the term $V_{12}$. We write the integral and apply the transformation (6.11). We get a sum of a $C^\omega$-smooth term and the integral

\[ J(u, v) = \int_{|z| \geq 2} \int_{\mathbb{R}} u_k(is) v_k(it) A\left(-\frac{1}{2}(k + is)\right)^{-1} A\left(-\frac{1}{2}(k + it)\right)^{-1} \times \left(\frac{z}{\bar{z}}\right)^k \left|\frac{z}{\bar{z}}\right|^{-2is-it} d\bar{z} \, dt \, ds. \]

As above, we change $s \mapsto s - i\varepsilon$ in the box and get integrals $J(u, v, \varepsilon)$ with $\varepsilon > 0$. As above,

\[ J(u, v; \varepsilon) = \int_{\mathbb{R}} u_k(is) v_k(it) A\left(-\frac{1}{2}(k + is)\right)^{-1} A\left(-\frac{1}{2}(k + it)\right)^{-1} \times \left(\frac{z}{\bar{z}}\right)^k \left|\frac{z}{\bar{z}}\right|^{-2is-it} d\bar{z} \, dt \, ds. \]

If $k > 0$, then the term in square brackets is zero (we pass to polar coordinates and get 0 after the integration with respect to the angle coordinate). If $k = 0$, then we get

\[ \frac{2^{-\varepsilon-is-it}}{\varepsilon + i(s + t)}. \]

However, supp$(u_0)$, supp$(v_0)$ are contained in domains $s > 0$, $t > 0$, and actually we have no singularity. Thus $V_{12}$ is $C^\omega$-smooth.

The same examination shows $C^\omega$-smoothness of $V_{21}$. This completes the proof of Lemma 6.4. \qed
7. Asymptotics of the kernel in the parameters

7.1. The statement. Let us modify a notation for the kernel $\mathcal{K}$. Set

$$
\mathcal{K}^\circ(z; \lambda; \sigma) :=
$$

$$
:= \frac{1}{\Gamma^C(a+b|a+b)} 2F_1\left[ a + \lambda + \frac{\sigma}{2} | a - \bar{\lambda} + \frac{\sigma}{2}, a - \lambda - \frac{\sigma}{2} | a + \bar{\lambda} - \frac{\sigma}{2}; z \right] =
$$

$$
= \frac{1}{\Gamma^C(a+b|a+b)} 2F_1\left[ a + \frac{k+\sigma+iz}{2} | a + \frac{-k+\sigma+iz}{2}, a + \frac{-k-\sigma-is}{2} | a + \frac{k-\sigma-is}{2}; z \right],
$$

where $\lambda \in \Lambda$, $\sigma \in \mathbb{R}$. In fact,

$$
\mathcal{K}^\circ(z; \frac{k+is}{2}; \sigma) = \mathcal{K}(z; k, \sigma + is).
$$

However, in calculations of this section the variables $\sigma$ and $is$ have different meanings.

Denote

$$
t_\pm(z) = 1 \pm \sqrt{1-1/z}.
$$

Theorem 7.1. Then for a fixed $z$ we have the following asymptotic expansion

$$
\mathcal{K}^\circ(z; \lambda; \sigma) = \frac{1}{\Gamma^C(a - \lambda - \frac{\sigma}{2} | a + \bar{\lambda} - \frac{\sigma}{2}) \Gamma^C(b + \lambda + \frac{\sigma}{2} | b - \bar{\lambda} + \frac{\sigma}{2}) \cdot |\lambda|^N} \times
$$

$$
\times \frac{|1 - 1/z|^{-1/2} \cdot |1 - z|^{b-a} \cdot |z|^{-a-b}}{1}
$$

$$
\times \left[ \frac{t_- (z)}{t_+ (z)} \right] ^{\frac{\sigma}{2} + \lambda - \frac{\sigma}{2} - \bar{\lambda}} \sum_{k \geq 0, l \geq 0, k+l \leq N} \frac{\lambda^{-k} \bar{\lambda}^{-l}}{k! l!} A_k(\sigma, \sqrt{1-z}) A_l(\sigma, \sqrt{1-z}) +
$$

$$
+ \left[ \frac{t_+ (z)}{t_- (z)} \right] ^{\frac{\sigma}{2} + \lambda - \frac{\sigma}{2} - \bar{\lambda}} \sum_{k \geq 0, l \geq 0, k+l \leq N} \frac{\bar{\lambda}^{-k} \lambda^{-l}}{k! l!} A_k(\sigma, -\sqrt{1-z}) A_l(\sigma, -\sqrt{1-z}) +
$$

$$
+ R_N(z, \sigma, \lambda),
$$

where $A_k(\xi)$ are rational expressions in $\xi$ (depending on the parameters $a, b$) having poles at $\xi = 0, \pm 1$ and $A_0 = 1$. The reminder $R_N(z)$ satisfies

$$
R_N(z, \sigma, \lambda) = O(|\lambda|^{-N}), \quad \text{as } \lambda \to \infty,
$$

moreover $O(\cdot)$ is uniform in $z$ and $\sigma$ on compact subsets in $\mathbb{C} \times \mathbb{R}$.

The proof occupies the rest of this section.

Remark. This formula is a counterpart of Watson’s [44] formula for asymptotics of the Gauss hypergeometric functions $2F_1[a - \lambda, b + \lambda; c; z]$ in the parameter $\lambda$ (see an exposition of Watson’s results in [29], Sect. 7.2, see also a remark in [41], p.162, on typos in [44]). We do not see a way to reduce our statement to Watson’s work.

$\Box$

Remark. Lemma 2.1 gives us an asymptotics of the gamma-factor in (7.1).

7.2. Stationary phase approximation. We transform $\mathcal{K}^\circ(z, \lambda, \sigma)$ as

$$
\int_{\mathbb{C}} R(t, z, \sigma) \exp\left\{ Q(t, z, \lambda, \sigma) \right\} dt,
$$

where

$$
R(t, z) := t^{a - \frac{\sigma}{2} - 1} |a - \frac{\sigma}{2} - 1| (1 - t)^{b + \frac{\sigma}{2} - 1} |b + \frac{\sigma}{2} - 1| (1 - tz)^{-a - \frac{\sigma}{2}}
$$

and

$$
Q(t, z, \lambda, \sigma) :=
$$
and

\[
(7.4) \quad Q(t, z, \lambda) := \lambda \ln \left( \frac{t(1 - zt)}{1 - t} \right) - \ln \left( \frac{t(1 - zt)}{1 - t} \right) = \lambda \ln \left( \frac{t(1 - zt)}{1 - t} \right) + i \frac{s}{2} \ln \left( \frac{t(1 - zt)}{1 - t} \right).
\]

The function \( \text{Im} \ln(\ldots) \) is ramified, however the exponent is well-defined and formulas below contain only partial derivatives of \( \ln(\ldots) \), which are independent of the choice of a branch.

We apply the stationary phase approximation, see, e.g., Fedoryuk [10], Hörmander [19]. Singular points are 0, 1, \( \infty \). Stationary points are

\[
t_{\pm} = 1 \pm \sqrt{1 - 1/z},
\]

they are the same for both summands in (7.4). This could be a fatal obstacle for an evaluation of a uniform asymptotics, however this does not happen. Also the domain of convergence of the integral (7.2) is smaller than it is necessary for our purposes.

Consider a partition of unity

\[
1 = \rho_0 + \rho_1 + \rho_\infty + \rho_+ + \rho_- + \tau,
\]

where \( \rho_\alpha \) is zero outside a small neighborhood of \( \alpha \), and \( \tau \) is zero in neighborhoods of 0, 1, \( z^{-1} \), \( \infty \), \( t_{\pm} \). According to this partition we expand (7.2) into a sum of 7 integrals,

\[
J = I_0 + I_1 + I_{z^{-1}} + I_\infty + I_+ + I_- + J.
\]

Obviously (see [10], Lemma III.2.1), for each \( N \) we have

\[
J = O(k^2 + s^2)^{-N} \quad \text{as } n + is \to \infty.
\]

### 7.3. Preparatory statement.

**Theorem 7.2.** Let \( \Omega \) be a domain in \( \mathbb{C} \), \( f(t) \), \( \varphi(t) \) be holomorphic in \( \Omega \). Let \( t_0 \) be a unique zero of \( \varphi'(t) \) in \( \Omega \) and \( \varphi''(t_0) \neq 0 \). Let \( \rho(t) \) be a \( C^\infty \)-smooth function compactly supported by \( \Omega \) such that \( \rho = 1 \) in a neighborhood of \( t_0 \). Consider the integral

\[
(7.5) \quad I(\lambda) = \int_{\Omega} \rho(t) f(t) \overline{f(t)} \exp \left\{ i \operatorname{Re}(\lambda \varphi(t)) \right\} dt,
\]

where \( \lambda \in \mathbb{C} \) is a parameter. Then

a) For \( |\lambda| > 1 \) we have the following expansion

\[
(7.6) \quad I(\lambda) = \frac{1}{|f''(t_0)||\lambda|} \exp \left\{ i \operatorname{Re}(\lambda \varphi(t_0)) \right\} \times
\]

\[
\times \left( \sum_{k \geq 0, j \geq 0, k + j < N} \frac{\lambda^{-k} \varphi^{-j}}{k! j!} a_k(f, \varphi) a_j(f, \overline{\varphi}) + R_N(\lambda) \right),
\]

where \( a_k \) are rational expressions

\[
a_k = a_k(\varphi(t_0), \varphi'(t_0), \ldots; f(t_0), f'(t_0), \ldots; \varphi''(t_0)^{-1})
\]

and \( a_0 = 1 \). The reminder \( R_N \) satisfies

\[
(7.7) \quad R_N(\lambda) = O(|\lambda|^N) \quad \text{as } \lambda \to \infty.
\]
b) The asymptotic expansion

\[ I(\lambda) \sim |\lambda|^{-1} \sum_{k \geq 0, t \geq 0} \frac{c_{kl}}{\lambda^k} \] 

as \( \lambda \to \infty \)
can be written as

\[ (7.8) \quad I(\lambda) \sim \frac{1}{|f''(t_0)||\lambda|} \exp\left\{ i \Re(\lambda \varphi(t_0)) \right\} \times \]

\[ \times \exp\left\{ \frac{i}{2\lambda} \frac{\partial^2}{\partial t^2} \left( f(t) \exp\left\{ \lambda(\varphi(t) - \varphi'(t_0) - \frac{1}{2} \varphi''(t_0)(t - t_0)^2) \right\} \right) \right\}_{t=t_0} \times \]

\[ \times \exp\left\{ \frac{i}{2\lambda} \frac{\partial^2}{\partial t^2} \left( f(t) \exp\left\{ \lambda(\varphi(t) - \varphi'(t_0) - \frac{1}{2} \varphi''(t_0)(t - t_0)^2) \right\} \right) \right\}_{t=t_0} . \]

\[ (7.9) \quad S(x, x_0) := S(x) - \frac{1}{2} (H(x_0)(x - x_0), (x - x_0)), \]

this expression is the part of the Taylor expansion of \( S(x) \) at \( x_0 \) starting cubic terms. Then the following expansion take holds:

\[ (7.10) \quad I(\sigma) = \left( \frac{2\pi}{\sigma} \right)^{n/2} |\det H(x_0)|^{-1/2} \exp\left\{ \frac{i\pi}{4} \sgn H(x_0) \right\} \times \]

\[ \times \left( \sum_{k=0}^{N-1} \frac{\sigma^{-k}}{k!} L^k (f(x) \exp\{i\sigma S(x, x_0)\}) \right|_{x=x_0} + \sigma^{-N+[2N/3]} V(\sigma) \right) , \]

where \( V(\sigma) \) is bounded.

Let us return to our integral \( I(\lambda) \). Without loss of generality, we can set \( t_0 = 0 \), \( \varphi''(t_0) = 1 \), i.e.,

\[ \varphi(t) = \frac{1}{2} t^2 + r(t), \quad \text{where} \ r(0) = r'(0) = r''(0) = 0 . \]

Set \( \lambda = se^{i\theta} \), \( s > 0 \). Set \( z = x + iy \), then

\[ \varphi(x, y) = \frac{1}{2}(x^2 - y^2 + 2ixy) + r(x, y) . \]
Thus we come to an oscillating integral in $s$ with the parameter $\theta$.

\[ I(s, \theta) = \int \rho(x, y)f(x, y)f(x, y) \exp \left\{ is(\cos \theta \text{ Re} \varphi(x, y) + \sin \theta \text{ Im} \varphi(x, y)) \right\} \, dx \, dy. \]

We wish to apply the general statement formulated above. The Hessian is given by

\[ H = 2 \begin{pmatrix} \cos \theta & \sin \theta \\ \sin \theta & -\cos \theta \end{pmatrix}, \quad H^{-1} = \frac{1}{2} \begin{pmatrix} \cos \theta & \sin \theta \\ \sin \theta & -\cos \theta \end{pmatrix}. \]

The signature is 0. The differential operator $L$ is

\[ L = \frac{i}{4} \left( \cos \theta \left( \frac{\partial^2}{\partial x^2} - \frac{\partial^2}{\partial y^2} \right) + 2 \sin \theta \frac{\partial^2}{\partial x \partial y} \right) = \frac{i}{2} \left( e^{i\theta} \frac{\partial^2}{\partial t^2} + e^{-i\theta} \frac{\partial^2}{\partial t} \right). \]

Next, we rewrite our phase function $S(\cdot)$ as

\[ e^{-i\theta} \varphi(t) + e^{i\theta} \overline{\varphi(t)}. \]

Therefore the expression (7.9) is

\[ e^{-i\theta} r(t) + e^{i\theta} \overline{r(t)}. \]

Applying (7.10), we get

\[
I(s, \theta) := \frac{2\pi}{s} \exp \left\{ \frac{i}{2s} \left( e^{i\theta} \frac{\partial^2}{\partial t^2} + e^{-i\theta} \frac{\partial^2}{\partial t} \right) \right\} \times
\times \left( f(t) \overline{f(t)} \exp \left\{ is(e^{-i\theta} r(t) + e^{i\theta} \overline{r(t)}) \right\} \right) \bigg|_{t=0} = \\
= \frac{2\pi}{s} \exp \left\{ \frac{i}{2se^{-i\theta}} \frac{\partial^2}{\partial t^2} \right\} \left( f(t) \exp \left\{ ise^{-i\theta} r(t) \right\} \right) \bigg|_{t=0} \times
\times \exp \left\{ \frac{i}{2se^{i\theta}} \frac{\partial^2}{\partial t} \right\} \left( \overline{f(t)} \exp \left\{ ise^{i\theta} \overline{r(t)} \right\} \right) \bigg|_{t=0}.
\]

We obtained asymptotics in $s$ for fixed $\theta$. However, $\theta$ ranges in a compact set, by [10], Theorem III.2.2, we get that the term $V(\cdot)$ in (7.10) is bounded uniformly in $\theta$.

a) follows from b).

c) We again refer to the parametric version of the stationary phase approximation, see [10], Theorem III.2.2. \(\square\)

7.4. Contribution of the stationary points. Let us apply Theorem 7.2 to our integral (7.2). We have

\[
f(t) = R(t, z) = \left( \frac{t}{1-zt} \right)^a (1-t)^b \left( \frac{1-t}{t(1-zt)} \right)^{\frac{a}{2}} (t(t-1))^{-1};
\]

\[
\varphi(t) = 2 \ln \left( \frac{1-t}{t(1-zt)} \right).
\]

Denote

\[
\zeta = \sqrt{1-1/z}.
\]
We substitute $t = t_+$ and transform the factors of $R(t, z) = f(t)\overline{f(t)}$:

\begin{equation}
\left.\frac{t}{1-zt}\right|_{t=t_+}^{a|a|} = \left(\frac{1-\zeta^2}{-\zeta}\right)^{a|a|} = ((z-1)\bar{z})^{-a/2-b/2};
\end{equation}

\begin{equation}
(1-t)^{b|b|}\bigg|_{t=t_+} = \left(\frac{1-\bar{z}}{z}\right)^{b/2|b/2};
\end{equation}

\begin{equation}
\left.\frac{1-t}{t(1-zt)}\right|_{t=t_+}^{z}\bigg|_{t=t_+} = \left(\frac{1-\zeta}{1+\zeta}\right)^{z/2} ;
\end{equation}

\begin{equation}
\left.\left((t-1)^{-1|-1}\right)\bigg|_{t=t_+} = \left(\frac{-1}{\zeta(1+\zeta)}\right)^{1|1}.\right.
\end{equation}

Next,

$$\varphi(t_+) = 2\ln\left(\frac{1-\zeta}{1+\zeta}\right),$$

therefore

$$\exp\left\{i\Re(\varphi(t_+))\frac{k}{2}(k+is)\right\} = \left(\frac{1-\zeta}{1+\zeta}\right)^{\lambda|-\lambda} = \left(\frac{t_+}{t_+}\right)^{\lambda|-\lambda}.$$ 

Finally,

$$\varphi''(t) = \frac{-2}{(1-t)^2} + \frac{2}{i^2} + \frac{2z^2}{(1-tz)^2},$$

and

$$\varphi''(t_+) = \frac{-4}{\zeta(1+\zeta)^2}.$$ 

Uniting these data we get that the leading term at the point $t_+$ is

\begin{equation}
-|\zeta| |1-z|^{b-a} |z|^{-a-b} \left(\frac{t_+}{t_+}\right)^{\lambda|-\lambda} \frac{1}{(k^2+s^2)^{1/2}}.
\end{equation}

The general form of the asymptotic expansion at $t = t_+$ follows from Theorem 7.2.

**7.5. Contributions of the singular points.**

**Lemma 7.3.** Contributions at the singular points $0$, $1$, $\infty$ are $O(|\lambda|^{-N})$ for any $N$.

**Proof.** For definiteness examine the point $0$. We have the integral

$$I_0(\lambda) = \int_C \rho_0(t) t^{a-1|a-1} (1-t)^{c-a-1} (1-zt)^{-a} \left(\frac{t(1-zt)}{1-t}\right)^{\lambda|-\lambda} d\bar{t},$$

defined as an analytic continuation. Keeping in mind that a support of $\rho_0$ can be chosen sufficiently small, we pass to a new variable in a neighborhood of $0$,

$$u = \frac{t(1-zt)}{1-t},$$

and come to an integral of the form

$$I_0(\lambda) = \int_C u^{a-\lambda-1|a+\lambda-1} \Phi(u) d\bar{u},$$

where $\Phi$ is a smooth compactly supported function. It remains to apply Theorem 7.2.

Argumentation for other singular points is the same. $\square$
8. Symmetry of difference operators

Here we prove Theorem 1.7, i.e., show that if \( f \in \mathcal{D}(\hat{C}) \), then \( J_{a,b}f \) is contained in the space \( W_{a,b} \) of meromorphic functions on \( \Lambda_C \). Also we show that \( \mathcal{E} \) and \( \mathcal{F} \) are formally adjoint one to another on \( W_{a,b} \), see Theorem 8.4.

8.1. Beginning of the proof of Theorem 1.7. We follow the list of properties in the definition of \( W_{a,b} \), see Subsect. 1.3.

a) is a corollary of the symmetry \( \mathcal{K}_{a,b}(z; -k, -\sigma) = \mathcal{K}_{a,b}(z; k, \sigma) \).

b) We must examine poles of \( \mathcal{K}_{a,b}(z; k, \sigma) \) as a function of the variable \( \sigma \) for a fixed \( z \in \hat{C} \), \( k \in \mathbb{Z} \). Let \( a + b \neq 1 \). We look to the expansion (3.26) of \( _2F_1[] \) at \( z = 0 \). The only source of poles of \( \mathcal{K} \) are zeros of the denominators in (3.28), i.e., zeros of the expression

\[
R(k, \sigma) := \Gamma_C(a + \frac{k + \sigma}{2}) \Gamma_C(a + \frac{k - \sigma}{2}) \times \Gamma_C(b + \frac{k + \sigma}{2}) \Gamma_C(b + \frac{k - \sigma}{2}).
\]

This gives us the desired list of possible poles.

Let us examine the case \( a + b = 1 \). The decomposition of the hypergeometric functions (3.26) at \( z = 0 \) produces an expression of the type

\[
\mathcal{K}_{a,b}(z; k, \sigma) = \frac{u_{a,b}(z, k, \sigma) - v_{a,b}(z, k, \sigma)}{a + b - 1}
\]

with \( u_{a,b}, v_{a,b} \) having poles at zeros of \( R(k, \sigma) \). A decomposition at \( z = 1 \) gives

\[
\mathcal{K}_{a,b}(z; k, \sigma) = \frac{U_{a,b}(z, k, \sigma) - V_{a,b}(z, k, \sigma)}{a - b},
\]

therefore the singularity in (3.22) at \( a + b = 1 \) is removable.

d) Indeed, we have \( \mathcal{K}_{a,b}(p, q) = \mathcal{K}_{a,b}(q, p) \), i.e.,

\[
\frac{2}{\Gamma_C} \left[ a + \frac{p + q}{2} a + \frac{p - q}{2} \; a + b \; z \right] = \frac{2}{\Gamma_C} \left[ a + \frac{p + q}{2} a + \frac{p - q}{2} \; a + b \; z \right].
\]

This is a special case of the symmetry (3.30).

We also mention the following similar identity for (8.1):

\[
R(p, q) = R(q, p),
\]

it is a special case of (3.37).

The statement c) about the behavior at infinity is a corollary of the expansion (7.1) and the following lemma

**Lemma 8.1.** Let \( t_\pm(z) \) be as in Theorem 7.1. Let \( \Phi \in \mathcal{D}(\hat{C}) \) be a function with a simply connected support. Then for any \( A > 0 \) for any \( N > 0 \) in the strip \( |\text{Re} \, \sigma| < A \) we have

\[
\int_{\hat{C}} \Phi(z) \left( \frac{t_-(z)}{t_+(z)} \right)^{(k+\sigma)/2} \, d\mathbb{F} = O\left( k^2 + (\text{Im} \, \sigma)^2 \right)^{-N}
\]
as \( (k^2 + (\text{Im} \, \sigma)^2) \to \infty \).
We need a simply connected support since the integrand is ramified at the points \( z = 0, \ z = 1 \). A proof of the lemma requires some preparations.

8.2. A change of variable. We define a new variable

\[
p := \frac{t_+(z)}{t_-(z)},
\]

The inverse map is done by

\[
(8.5) \quad z = \zeta(p) := \frac{(p+1)^2}{4p}.
\]

The map \( \zeta(p) \) determines a two-sheet covering map from

\[
(8.6) \quad \hat{\mathbb{C}} := \mathbb{C} \setminus \{0, 1, -1\}
\]

to \( \hat{\mathbb{C}} \). Notice that

\[
1 - z = -\frac{(p-1)^2}{4p}, \quad \sqrt{1 - \frac{1}{z}} = \frac{p - 1}{p + 1}, \quad \zeta'(p) = \frac{p^2 - 1}{4p^2},
\]

\[
(8.7) \quad t_+ = \frac{2p}{p + 1}, \quad t_- = \frac{2}{p + 1}, \quad \frac{t_+}{t_-} = p.
\]

Also,

\[
(8.9) \quad \zeta(p^{-1}) = \zeta(p), \quad \zeta'(p^{-1}) p^{-1} = \zeta(p) p.
\]

8.3. Proof of Theorem 1.7.c.

Proof of Lemma 8.1. We substitute \( z = \zeta(p) \) to the integral and get

\[
\frac{1}{16} \int_{\hat{\mathbb{C}}} p^{(k+\sigma)/2} \left| \Phi(\zeta(p)) \right| p^2 - 1 |p^{-2} |\, dp.
\]

This is a Mellin transform of a function compactly supported by \( \hat{\mathbb{C}} \). In virtue of Theorem 2.2, the integral rapidly decreases in the union of strips \( | \text{Re} \sigma | < A \).

Proof of the statement c) of Theorem 1.7. We represent \( \varphi(z) \) as a sum of functions in \( \mathcal{D}(\hat{\mathbb{C}}) \) with simply connected supports. Next, we decompose the kernel according to Theorem 7.1 and apply the lemma to each summand.

8.4. Continuity.

Corollary 8.2. The map \( J_{a,b} \) is a continuous map from \( \mathcal{D}(\hat{\mathbb{C}}) \) to \( L^2_{\text{even}}(\Lambda, \kappa_{a,b}) \).

Proof. Define the following seminorms on the space of smooth functions on \( \Lambda \):

\[
p_{a,N}(F) = \sup_{\lambda \in \Lambda} \left| \frac{\partial^N F}{\partial \sigma^N} (1 + |\lambda|)^a \right|,
\]

and the space \( Y \) defined by these seminorms. Clearly, our proof provides a continuity of \( J_{a,b} \) as a map \( \mathcal{D}(\hat{\mathbb{C}}) \) to \( Y \). It remains to notice that the identical embedding \( f \mapsto f \) of \( Y \) to \( L^2 \) is continuous.

If \( k = 0 \) and \( a = 1 \) (or \( b = 1 \)), then elements of \( \mathcal{W}_{a,b} \) have a pole of order two at \( k = 0, \ s = 0 \). In this case we write \( \lambda^2 F \) instead of \( F \) in the definition of the seminorms.

\[ \text{At the same point the spectral density has a zero of order 4.} \]
8.5. Invariance of $\mathcal{W}_{a,b}$. Consider the difference operators $\mathfrak{L}, \overline{\mathfrak{L}}$ defined above (1.28),

\begin{align}
\mathfrak{L} F(k,\sigma) &= \frac{(a + \frac{k+\sigma}{2})(b + \frac{k+\sigma}{2})}{(k+\sigma)(1+k+\sigma)}(F(k+1,\sigma + 1) - F(k,\sigma)) + \\
&\quad + \frac{(a + \frac{-k-\sigma}{2})(b + \frac{-k-\sigma}{2})}{(-k-\sigma)(1-k-\sigma)}(F(k-1,\sigma - 1) - F(k,\sigma));
\end{align}

\begin{align}
\overline{\mathfrak{L}} F(k,\sigma) &= \frac{(a + \frac{-k+\sigma}{2})(b + \frac{-k+\sigma}{2})}{(-k+\sigma)(1-k+\sigma)}(F(k-1,\sigma + 1) - F(k,\sigma)) + \\
&\quad + \frac{(a + \frac{k-\sigma}{2})(b + \frac{k-\sigma}{2})}{(k-\sigma)(1+k-\sigma)}(F(k+1,\sigma - 1) - F(k,\sigma)).
\end{align}

**Lemma 8.3.** The space $\mathcal{W}_{a,b}$ is invariant with respect to the operators $\mathfrak{L}, \overline{\mathfrak{L}}$.

**Proof.** Since $F(0,-1) = F(1,0) = F(-1,0) = F(0,1)$, the expressions

\[
\frac{F(k+1,\sigma + 1) - F(k,\sigma)}{1+k+\sigma}, \quad \frac{F(k-1,\sigma - 1) - F(k,\sigma)}{1-k-\sigma}
\]

have no poles at $k = -1, \sigma = 0$ and $k = 1, \sigma = 0$ respectively.

Since a function $F(k,\sigma)$ is even, it can not have a pole of order 1 at $k = 0, \sigma = 0$.

New poles of $F(k+1,\sigma + 1)$ that are not poles of $F(k,\sigma)$ are annihilated by the rational factor in (8.10).

The condition $\mathfrak{L} F(p,q) = \mathfrak{L} F(q,p)$ follows from a straightforward calculation. \(\square\)

8.6. Symmetry.

**Theorem 8.4.** For $(a,b) \in \Pi$, for $F, G \in \mathcal{W}_{a,b}$ we have

\begin{equation}
\langle \mathfrak{L} F, G \rangle_{L^2(\Lambda_c, d\mathcal{K}_{a,b})} = \langle F, \overline{\mathfrak{L}} G \rangle_{L^2(\Lambda, d\mathcal{K}_{a,b})}.
\end{equation}

**Corollary 8.5.** Operators $\frac{1}{4}(\mathfrak{L} + \overline{\mathfrak{L}})$, $\frac{1}{4}(\mathfrak{L} - \overline{\mathfrak{L}})$ are symmetric on the $J_{a,b}$-image of $\tilde{D}(\mathbb{C})$.

**Remark.** In fact, the proof uses only properties of $F \in \mathcal{W}_{a,b}$ in strips $|\text{Re}\, \sigma| < 1 + \varepsilon$. So we can define operators $\mathfrak{L}, \overline{\mathfrak{L}}$ on a space of meromorphic functions in the strip satisfying an obvious list of conditions. \(\Box\)

8.7. **Proof of Theorem 8.4 for the case $(a,b) \in \Pi_{\text{cont}}**. First, we notice that for pure imaginary $\sigma$ we have $G(k,\sigma) = G(k,-\sigma)$, the last function is meromorphic and also is contained in $\mathcal{W}_{a,b}$. Let $R(k,\sigma)$ be given by (8.1). Then

\begin{align}
4\pi^2i \langle \mathfrak{L} F, G \rangle &= \sum_k \int_{\mathbb{R}} \left\{ \frac{(a + \frac{k+\sigma}{2})(b + \frac{k+\sigma}{2})}{(k+\sigma)(1+k+\sigma)}(F(k+1,\sigma + 1) - F(k,\sigma)) + \\
\quad + \frac{(a + \frac{-k+\sigma}{2})(b + \frac{-k+\sigma}{2})}{(-k+\sigma)(1-k+\sigma)}(F(k-1,\sigma - 1) - F(k,\sigma)) \right\} \\
&\quad \times G(k,-\sigma) (k-\sigma)(k+\sigma) R(k,\sigma) d\sigma.
\end{align}

Let us expand the expression in the curly brackets $\{\ldots\}$ as a sum of 4 summands that include $F(k+1,\sigma + 1)$, $F(k,\sigma)$, $F(k-1,\sigma - 1)$, $F(k,\sigma)$. The whole expression $\{\ldots\}$ is holomorphic near the contour of integration. The summands have simple poles on the contour, and we pass to an integration in the sense of principal values.
Let us examine the summand corresponding $F(k + 1, \sigma + 1)$. We get

\begin{equation}
\sum_k \text{v.p.} \int_{i \eta} \frac{k - \sigma}{1 + k + \sigma} F(k + 1, \sigma + 1) \tilde{G}(k, -\sigma) \tilde{R}(k, \sigma) \, d\sigma,
\end{equation}

where

\begin{equation}
\tilde{R}(k, \sigma) := \left( a + 1 + \frac{k + \sigma}{2} \right) \left( b + 1 + \frac{k + \sigma}{2} \right) R(k, \sigma) =
\end{equation}

\begin{equation}
= \Gamma^C(a + 1 + \frac{k + \sigma}{2}) \Gamma^C(a + \frac{k - \sigma}{2}) \Gamma^C(b + 1 + \frac{k + \sigma}{2}) \Gamma^C(b + \frac{k - \sigma}{2}).
\end{equation}

**Lemma 8.6.** For $0 < a < 1$, $0 < b < 1$ the integrand in (8.14) has no poles in the strip $-1 < \text{Im} \sigma < 0$.

**Proof.** We enumerate possible (simple) poles of the factors.

a) Factor $\tilde{G}(k, -\sigma)$. In this case we can have poles if $k = 0$. Since $a < 1$, $b < 1$ the poles $2 - 2a$, $2 - 2b$ are outside our strip. On other hand the pole $2a - 2$ (resp. $2b - 2$) is contained in the strip if $1/2 < a < 1$ (resp. if $1/2 < b < 1$).

b) Factor $F(k + 1, \sigma + 1)$ has a pole in our strip for $k = -1$ at $\sigma = 2a - 1$ (resp. $\sigma = 2b - 1$) if $0 < a < 1/2$, (resp. $0 < b < 1/2$).

c) Since $a > 0$, $b > 0$ the expression $\tilde{R}(k, \sigma)$ has no poles in our strip.

However, the poles of $\tilde{G}(k, -\sigma)$ and of $F(k + 1, \sigma + 1)$ are zeros of $\tilde{R}(k, \sigma)$. Therefore the product is holomorphic. □

**Lemma 8.7.** In (8.14), we can change the integration contour to $1 + i \mathbb{R}$.

**Proof.** The integrand has no poles between contours $i \eta$ and $1 + i \eta$, but has poles on contours, the integral is taken in the sense of principal values. We have only two such poles, $\sigma = 0$ on the contour $i \eta$ for $k = -1$ and $\sigma = -1$ for $k = 0$.

Thus the difference between the two integrals is $2\pi$ by half of the sum of residues, i.e.,

\begin{equation}
\frac{2\pi}{2} \left\{ (-1 - \sigma) F(0, \sigma + 1) \tilde{G}(-1, -\sigma) \left( a + 1 + \frac{-1 - \sigma}{2} \right) \left( b + 1 + \frac{-1 - \sigma}{2} \right) R(-1, \sigma) \bigg|_{\sigma = 0} +
\right.
\end{equation}

\begin{equation}
+ (0 - \sigma) F(1, \sigma + 1) \tilde{G}(0, -\sigma) \left( a + 1 + \frac{0 - \sigma}{2} \right) \left( b + 1 + \frac{0 - \sigma}{2} \right) R(0, \sigma) \bigg|_{\sigma = -1} \right\}.
\end{equation}

Let us show that the sum is zero. Since $F$, $G$ are even and satisfy (12.23), we have

\begin{equation}
F(0, 1) = F(1, 0), \quad \tilde{G}(-1, 0) = \tilde{G}(0, -1).
\end{equation}

By (8.23), we have

\begin{equation}
R(-1, 0) = R(0, -1).
\end{equation}

The remaining factors give

\begin{equation}
-(a + \frac{1}{2})(b + \frac{1}{2}) \quad \text{and} \quad (a + \frac{1}{2})(b + \frac{1}{2}),
\end{equation}

i.e., the same expressions with different signs. □

**End of the proof of Theorem 8.2** for the case $(a, b) \in \Pi_{\text{cont}}$. Thus we can replace the integration in (8.14) by the integration over the contour $-1 + i \mathbb{R}$. 
We change the variables $l = k + 1, t = \sigma + 1$ and get
\[
\sum_{\ell} \text{v.p.} \int_{i\mathbb{R}} \frac{l - t}{-1 + l + t} F(l, t) G(l - 1, -t + 1) \tilde{R}(l - 1, t - 1) dt.
\]
Next,
\[
\tilde{R}(l - 1, t - 1) = R(l, t) (a + \frac{\gamma}{2} - \frac{l - t}{2}) (b + \frac{\gamma}{2} - \frac{l - t}{2}),
\]
and we come to
\[
\sum_{\ell} \text{v.p.} \int_{i\mathbb{R}} F(l, t) \left[ \frac{(a + \frac{\gamma}{2} - \frac{l - t}{2}) (b + \frac{\gamma}{2} - \frac{l - t}{2})}{(-l - t)(1 - l - t)} G(l - 1, -t + 1) \right] \times
\]
\[
\times (l - t)(l + t) R(l, t) dt.
\]
We transform the expression in the big brackets to the form $U(l, -t)$, where
\[
U(l, t) = \frac{(a + \frac{\gamma}{2} - \frac{l + t}{2}) (b + \frac{\gamma}{2} - \frac{l + t}{2})}{(-l + t)(1 - l + t)} G(l - 1, t + 1).
\]
Thus we finished the transformation of the summand of the (8.13) corresponding to $F(k, \sigma)$. The transformation of the summand corresponding to $F(k - 1, \sigma - 1)$ is similar. The case of the summands $F(k, \sigma)$ is obvious. We come to the desired expression.

8.8. **End of the proof of Theorem 8.4.** Due to the homographic transformations, it is sufficient to examine the case $a < 0$. Let $\Phi, \Psi \in \mathcal{W}_{a,b}$. Denote
\[
(8.16) \quad U(a, b; k, \sigma) := \Phi(k, \sigma) \Psi(k, -\sigma) \varphi_{a,b}(k, \sigma).
\]
For $(a, b) \in \Pi_{\text{cont}}$ we have
\[
(8.17) \quad 4\pi^2 i \langle \Phi, \Psi \rangle_{L^2(\Lambda, \varphi_{a,b})} = \sum_{k} \int_{i\mathbb{R}} U(a, b; k, \sigma) d\sigma.
\]
We wish to write the analytic continuation of this expression to the domain $(a, b) \in \Pi, a \leq 0$.

Possible singularities of $U$ as a function in $\sigma$ in the strip $|\text{Re} \sigma| < 1$ are the following:

- If $b > 1/2$, then both functions $\Phi, \Psi$ have poles at $(k, \sigma) = (0, \pm(2 - 2b));$
- $\varphi_{a,b}(k, \sigma)$ has poles at $(k, \sigma) = (0, \pm 2a)$.

Due to our restrictions $2b - 2 < 2a < -2a < 2 - 2b$.

Thus all summands of (8.17) except 0-th are holomorphic in $|a| < 1 - b$.

**Lemma 8.8.** Fix $b$. Assume that $\Phi, \Psi$ be even rapidly decreasing meromorphic functions in the strip $|\text{Re} \sigma| < 1$ satisfying the condition (1.25) and having poles only at the points $(0, \pm(2 - 2b))$. Then the following expression is holomorphic in the domain $|a| < 1 - b$:
\[
(8.18) \quad \gamma^b(a) := \begin{cases} 
\int_{i\mathbb{R}} U(a, b; 0, \sigma) d\sigma, & \text{if } a \geq 0, \\
\int_{i\mathbb{R}} U(a, b; 0, \sigma) d\sigma + 4\pi i \text{res}_{\sigma = -2a} U(a, b; 0, \sigma), & \text{if } a \leq 0.
\end{cases}
\]

**Proof of Lemma 8.8.** Denote
\[
\gamma_{i\mathbb{R}}(a) = \int_{i\mathbb{R}} U(a, b; 0, \sigma) d\sigma, \quad \Xi_{\pm}(a) := 2\pi i \text{res}_{\sigma = \pm 2a} U(a, b; 0, \sigma).
\]
Since $U$ is even in $\sigma$, we have $\Xi_-(a) = -\Xi_+(a)$. Due to the factor $(k + \sigma)(k - \sigma)$ in the Plancherel density, we have $\Xi_+(0) = 0$. Therefore $\Xi_{\pm}(a)$ are holomorphic in the disk $|a| < 1 - b$.

Consider a contour $L$ on the plane $\sigma \in \mathbb{C}$ composed of the ray $(-\infty, b - 1 + \varepsilon]$, the upper half of the circle $|\sigma| = 1 - b - \varepsilon$ and the ray $[1 - b - \varepsilon, +\infty]$. The function $\gamma_L(a) := \int_L U(a, b; 0, \sigma) d\sigma$ is holomorphic in $a$ for $|a| < 1 - b$. For $\text{Re} \ a > 0$ we have $\gamma_L(a) = \gamma_{iR}(a) - \Xi_+(a)$. For $\text{Re} \ a < 0$ we have $\gamma_L(a) = \gamma_{iR}(a) - \Xi_-(a)$. This gives us the analytic continuation. □

9. The operator $J_{a,b}$ is an isometry

Here we prove the second part of Theorem 1.3

9.1. Statement.

Lemma 9.1. Let $f, g$ be smooth compactly supported functions on $\hat{\mathcal{C}}$. Then

$$(J_{a,b} f, J_{a,b} g)_{L^2(\Lambda, \kappa_{a,b})} = (f, g)_{L^2(\hat{\mathcal{C}}, \mu_{a,b})}.$$ 

Here a way of a proof is simpler than in Section 6. We show that $J_{a,b}$ is a perturbation of a version of the Mellin transform.

9.2. Orthogonality of packets.

Lemma 9.2. Let $f, g \in D(\hat{\mathcal{C}})$. Let $\text{supp}(f) \cap \text{supp}(g) = \emptyset$. Then

$$(J_{a,b} f, J_{a,b} g)_{L^2(\Lambda, dK_{a,b})} = 0.$$ 

Proof. By Corollary 8.2 the operator $J_{a,b}$ is continuous as an operator $D(\hat{\mathcal{C}}) \to L^2(\Lambda, dK_{a,b})$, by Theorem 8.4 it is symmetric on the image of $D(\hat{\mathcal{C}}, \mu_{a,b})$. We consider the difference operators

and repeat the proof of Lemma 6.3. □

9.3. Decomposition of the kernel. Starting from this place we examine the restriction of $J_{a,b}f$ to $\Lambda$. Recall that the operator $J_{a,b}$ is defined by the formula

$$(9.1) \quad J_{a,b}f(\lambda) = \int_{\mathcal{C}} f(z) \mathcal{K}(z, \lambda) \mu_{a,b}(z) \, d\bar{\mu}.$$ 

Decompose the kernel $\mathcal{K}(z, \lambda)$ according to (7.1) with $N = 3$. We consider $\lambda \in \Lambda$, and therefore we set $\sigma = 0$. Denote by $\omega$ the factor depending on $\lambda$ in the front of the expansion. We have

$$(9.2) \quad \omega(\lambda) \omega(\Lambda) = \kappa_{a,b}^{-1}(\lambda).$$
Notice also that the expression in brackets in (7.1) has a singularity at \( \lambda = 0 \). Denote by \( \Theta(\lambda) \) a smooth function, which equals 0 for \( |\lambda| \leq 1/3 \) and 1 for \( |\lambda| \geq 1/2 \). Represent the kernel as

\[
\mathcal{K}(z, \lambda) = \omega(\lambda)|1 - z|^{b-a}|z|^{-a-b} \times
\]

\[
\times \left\{ \left[ \left( \frac{t_+(z)}{t_-(z)} \right)^{\lambda - \overline{\lambda}} + \left( \frac{t_-(z)}{t_+(z)} \right)^{\lambda - \overline{\lambda}} \right] + \Theta(\lambda) \left[ \left( \frac{t_+(z)}{t_-(z)} \right)^{\lambda - \overline{\lambda}} \sum_{k \geq 0, l \geq 0, 1 \leq k + l \leq 2} \frac{\overline{\lambda}^{-k}}{k!} \lambda^{-l} A_k(\sqrt{1 - z}) A_l(\sqrt{1 - \overline{z}}) + \left( \frac{t_-(z)}{t_+(z)} \right)^{\lambda - \overline{\lambda}} \sum_{k \geq 0, l \geq 0, 1 \leq k + l \leq 2} \frac{\overline{\lambda}^{-k}}{k!} \lambda^{-l} A_k(\sqrt{1 - \overline{z}}) A_l(\sqrt{1 - z}) + R_3(z, \lambda) \right] \right\},
\]

where \( R_3(z, \lambda) \) is a smooth function in \( z \in \mathcal{C} \) and \( \lambda \),

\[
R_3(z, \lambda) = O(|\lambda|^{-3}) \quad \text{as } \lambda \to \infty.
\]

uniformly on compact subsets of \( \mathcal{C} \). The summands corresponding to \( k = 0, l = 0 \) are smooth at \( \lambda = 0 \), so we do not multiply them by the patch function \( \Theta(\lambda) \).

Next, we change the variable as in (8.4)–(8.9):

\[
\zeta(p) = \frac{(p + 1)^2}{4p}
\]

and represent the operator \( J_{a,b} \) in the form

\[
J_{a,b} f(\lambda) = \omega(\lambda) \int_{\mathcal{C}} f(\zeta(p)) |1 - \zeta(p)|^{a-b-1/2} |\zeta(p)|^{-1/2} |\zeta'(p)|^2 \times
\]

\[
\times \left\{ \left[ \left( \frac{p^{1-\overline{\lambda}}}{p^{1-\lambda}} \right)^{\lambda - \overline{\lambda}} + \left( \frac{p^{1-\overline{\lambda}}}{p^{1-\lambda}} \right)^{\lambda - \overline{\lambda}} \right] + \Theta(\lambda) \left[ \left( \frac{p^{1-\overline{\lambda}}}{p^{1-\lambda}} \right)^{\lambda - \overline{\lambda}} \sum_{k \geq 0, l \geq 0, 1 \leq k + l \leq 2} \frac{\overline{\lambda}^{-k}}{k!} \lambda^{-l} A_k(\frac{p-1}{p+1}) A_l(\frac{p-1}{p+1}) + \left( \frac{p^{1-\overline{\lambda}}}{p^{1-\lambda}} \right)^{\lambda - \overline{\lambda}} \sum_{k \geq 0, l \geq 0, 1 \leq k + l \leq 2} \frac{\overline{\lambda}^{-k}}{k!} \lambda^{-l} A_k\left(\frac{p-1}{p+1}\right) A_l\left(\frac{p-1}{p+1}\right) \right] + R(\zeta(p), \lambda) \right\} d\overline{p},
\]

where \( \mathcal{C} \) denotes \( \mathbb{C} \setminus \{0, 1, -1\} \) as above.

It is convenient to split the operator \( J_{a,b} \) into a sum of operators,

(9.3) \[
J_{a,b} = [V_{0,0}^+ + V_{0,0}^-] + \sum V_{k,l}^+ + \sum V_{k,l}^- + V_{\text{rem}},
\]

where the summands correspond to the summands of the previous formula. We also denote

\[
\gamma(p) := |1 - \zeta(p)|^{a-b-1/2} |\zeta(p)|^{-1/2} |\zeta'(p)|^2.
\]

9.4. The main term.

**Lemma 9.3.** The operator \( \frac{1}{2\pi i} (V_{0,0}^+ + V_{0,0}^-) \) is a unitary operator from \( L^2(\mathbb{C}, \mu_{a,b}) \) to \( L^2_{\text{even}}(\Lambda, \alpha_{a,b}) \).
Therefore both integrals over $\tilde{\varphi}$ and we can apply the Plancherel formula for the Mellin transform. We come to

\begin{equation}
\left\langle (V^+_{0,0} + V^-_{0,0})f, (V^+_{0,0} + V^-_{0,0})g \right\rangle_{L^2(\Lambda, \omega)} = \int \left( \int_{\mathcal{C}} f(\zeta(p)) \frac{\gamma(p)(p^{\lambda-1} - \overline{p^{\lambda-1}} + p^{\lambda-1} \overline{p^{\lambda-1}})}{d\overline{p}} \right) \left( \int_{\mathcal{C}} \overline{g(\zeta(q)) g(q)(q^{-\lambda} + \overline{q^{-\lambda}} - \overline{p^{\lambda-1}})} \frac{d\overline{q}}{d\overline{p}} \right) d\lambda \tag{9.4}
\end{equation}

(we also applied (9.2). Transform this expression as

\begin{equation}
\int \left( \int_{\mathcal{C}} f(\zeta(p)) \frac{\gamma(p)|p|^2 (p^{\lambda-1} - \overline{p^{\lambda-1}} + p^{\lambda-1} \overline{p^{\lambda-1}})}{d\overline{p}} \right) \times \left( \int_{\mathcal{C}} \overline{g(\zeta(q)) g(q)|q|^2 (q^{-\lambda} + \overline{q^{-\lambda}} - \overline{p^{\lambda-1}})} \frac{d\overline{q}}{d\overline{p}} \right) d\lambda.
\end{equation}

Now we apply the remark about Mellin transforms of even functions from Subsect. 2.3. Keeping in mind (8.9), we observe that functions $f(\zeta(p))\gamma(p)|p|^2$ are $\times$-even. Therefore both integrals over $\mathcal{C}$ in (9.4) are Mellin transforms of even functions, and we can apply the Plancherel formula for the Mellin transform. We come to

\begin{equation}
\int_{\mathcal{C}} f(\zeta(p)) \overline{g(\zeta(p))} |\gamma(p)|^2 |p|^4 \frac{d\overline{p}}{|p|^2} = \int_{\mathcal{C}} f(\zeta(p)) \frac{1 - \zeta(p)^{2a-2b} |\zeta(p)|^{2a+2b-2} |\zeta'(p)|^2}{|1 - \zeta(p)|^{-1} |\zeta(p)|^{-1} |p|^2 |\zeta'(p)|^2} \frac{d\overline{p}}{d\overline{q}}.
\end{equation}

By (8.7) - (8.8) the expression in the big brackets is 1. Now we return to the variable $z = \zeta(p)$ and get the desired expression

\begin{equation}
\int_{\mathcal{C}} f(z)g(z) \mu_{a,b}(z) \frac{d\overline{p}}{d\overline{q}}.
\end{equation}

9.5. Other terms.

Lemma 9.4. The Hermitian form

\begin{equation}
\{f, g\} := \left\langle J_{a,b}f, J_{a,b}g \right\rangle_{L^2(\Lambda, \omega_{a,b})} - \left\langle (V^+_{0,0} + V^-_{0,0})f, (V^+_{0,0} + V^-_{0,0})g \right\rangle_{L^2(\Lambda, \omega_{a,b})}
\end{equation}

on $\mathcal{D}(\mathcal{C})$ can be written as

\begin{equation}
\{f, g\} = \int_{\mathcal{C}} \int_{\mathcal{C}} K(p, q) f(\zeta(p)) \overline{g(\zeta(q))} \frac{d\overline{p}}{d\overline{q}},
\end{equation}

where $K$ is a locally integrable function on $\mathcal{C} \times \mathcal{C}$ smooth outside the sets $p = q$, $p = q^{-1}$.

Proof. Expanding $J_{a,b}$ according to (9.3), we get many summands in (9.5). We wish to show that each summand can be written as (9.6) with its own $K$. Let
us start the discussion with the summand

\begin{equation}
(9.7)
\langle V_{0,0}^+, V_{0,1}^- \rangle_{L^2(\Lambda, \nu_\alpha, b)} = \int_\Lambda \left( \frac{1 - \Theta(\lambda)}{\lambda} \right) \left( \int_\mathbb{C} f(\zeta(p)) \gamma(p) \left| p^\lambda - 1 \right|^{-1} \frac{d\overline{p}}{\overline{p}} \right) \times \left( \int_\mathbb{C} \frac{g(\zeta(q)) \gamma(q) \left| q^{\lambda - 1} \right|^{-1}}{\left| q^{1/2} \right| + 1} \frac{d\overline{q}}{\overline{q}} \right) \lambda.
\end{equation}

The integral in the first big bracket is the Mellin transform of the function

\( F(p) := f(\zeta(p)) \gamma(p) |p|^2. \)

The integral in the second big bracket is a complex conjugate to the Mellin transform of

\( G(q) = g(\zeta(q^{-1})) \gamma(q^{-1}) |q|^{-2} A_1 \left( -\frac{1}{q} - 1 \right). \)

Thus we get

\begin{equation}
\langle V_{0,0}^+, V_{0,1}^- \rangle_{L^2(\Lambda, \nu_\alpha, b)} = \int_\Lambda \mathcal{M} F(\lambda) \overline{\mathcal{M} G(\lambda)} \frac{1 - \Theta(\lambda)}{\lambda} \lambda d\lambda.
\end{equation}

Denote by \( L(p) \) the inverse Mellin transform of \( \frac{1 - \Theta(\lambda)}{\lambda} \). It is easy to see that \( L(p) \) is an integrable function with a unique singularity of the type \( 1/(1 - p) \) at \( p = 1 \). We rewrite our integral as

\( \int_C \int_C L(pq) F(p) \overline{G(q)} \frac{d\overline{p}}{\overline{p}} \frac{d\overline{q}}{\overline{q}}, \)

and it has the desired form.

For other pairs \( V_{k,l}^\varepsilon, V_{k',l'}^{\varepsilon'} \), where \( \varepsilon, \varepsilon' = \pm 1 \), we have similar calculations. Instead of the boxed factor in (9.7), we get

\begin{equation}
(9.8)
\frac{1 - \Theta(\lambda)}{\lambda^{k + k'}}.
\end{equation}

For \( k + l + k' + l' \leq 2 \) we repeat the same considerations, in these cases inverse Mellin transforms of the functions (9.8) have (integrable) singularities\(^23\) at \( p = 1 \) of types

\( \frac{1}{1 - p}, \quad \frac{1}{1 - \overline{p}}, \quad \ln |1 - p|, \quad \frac{1}{1 - \overline{p}}. \)

If \( k + l + k' + l' \geq 3 \), then this expression is integrable in \( \lambda \), the triple integral is convergent, we can change the order of integrations and we immediately get an expression of the form (9.6) with real analytic \( K(p, q) \).

For the pairs including \( V_{\text{ren}} \), we get absolutely convergent triple integrals and analytic kernels \( K(p, q) \).

**9.6. Proof of Lemma 9.1** Now let \( f, g \in \mathcal{D}(\mathbb{C}) \) have disjoint supports. Then both terms in (9.5) are zero (see Lemma 9.2). Therefore the kernel \( K(p, q) \) satisfy the following property:

\( \{ f, g \} = \int_\mathbb{C} \int_\mathbb{C} K(p, q) \varphi(p) \overline{\psi(q)} dp dq = 0 \)

\(^23\)We can refer to corresponding formulas for the Fourier transform, see [14], Addendum, Sect. 1.7 (Russian edition) or [15], Sect. B.1.3 (English translation).
if $\varphi, \psi$ are $x$-even elements $\mathcal{D}(\bar{C})$ with disjoint supports.

We claim that $\{f, g\} = 0$ for any $x$-even functions $f, g \in \mathcal{D}(\bar{C})$. To observe this, we take a $x$-even partition of unity $\tau_j$ with small supports, and decompose

$$\{f, g\} = \sum_{k,l} \{\tau_k f, \tau_l g\}.$$ 

Clearly, we can make this sum as close to zero as we want by refinement of a partition of unity. We omit trivial details.

10. **Domains of self-adjointness**

Thus $J_{a,b}$ is unitary. Clearly the multiplication operators

$$f(z) \mapsto \frac{1}{2}(z + \bar{z})f(z), \quad f(z) \mapsto \frac{1}{2i}(z - \bar{z})f(z)$$

defined on $\mathcal{D}(\bar{C})$ are essentially self-adjoint in $L^2(C, \mu_{a,b})$ and commute. Therefore the operators $\frac{1}{2}(\mathcal{E} + \bar{\mathcal{E}})$, $\frac{1}{2i}(\mathcal{E} - \bar{\mathcal{E}})$ are essentially self-adjoint and commute on the subspace $J_{a,b}\mathcal{D}(\bar{C}) \subset L^2_{\text{even}}(\Lambda_C, dK_{a,b})$. But $\mathcal{W}_{a,b}$ contains this image. This establishes Theorem 1.2.a.

Theorem 1.8.a follows from the same argumentation.

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