Schrodinger equation as recurrences. I. Band-matrix Hamiltonians

This article has been downloaded from IOPscience. Please scroll down to see the full text article.
1983 J. Phys. A: Math. Gen. 164001
(http://iopscience.iop.org/0305-4470/16/17/016)
View the table of contents for this issue, or go to the journal homepage for more

Download details:
IP Address: 129.252.86.83
The article was downloaded on 30/05/2010 at 16:55

Please note that terms and conditions apply.

# Schrödinger equation as recurrences: I. Band-matrix Hamiltonians 

M Znojil<br>Institute of Nuclear Physics, Czechoslovak Academy of Sciences, 25068 Řež, Czechoslovakia

Received 31 January 1983, in final form 5 July 1983


#### Abstract

We describe a new algebraic treatment of the chain-model eigenvalue problem. It is based on the ( $2 t+1$ )-diagonal and asymptotically smooth structure of the corresponding Hamiltonians. The formalism is illustrated on the phenomenological doubly anharmonic oscillators.


## 1. Introduction and summary

One of the most common formal difficulties connected with the applicability of perturbation theory is usually the lack of a suitable unperturbed Hamiltonian $H_{0}$. A textbook example is the anharmonic oscillator $H_{\lambda}=p^{2}+r^{2}+\lambda r^{4}$ where the obvious choice $H_{0}=p^{2}+r^{2}$ leads to divergences (Simon 1969).

The puzzling divergence of the anharmonic-oscillator energies $E_{\lambda}$ may be removed easily by another choice of $H_{0}$ (Halliday and Suranyi 1980). Nevertheless, a more efficient solution of the Schrödinger eigenvalue problem $H_{\lambda} \psi=E_{\lambda} \psi$ was found by Graffi and Grecchi (1975). It employs directly the pentadiagonal structure of $H_{\lambda}$ in the unperturbed oscillator basis, and represents the solution in terms of the matrix continued fractions (MCF).

In practice the MCF recurrent technique is formally applicable to any ( $2 t+$ 1)-diagonal Hamiltonian, i.e. to the so-called chain models appearing in various branches of physics (see e.g. Haydock (1980) and references therein). A simple testing example of such a system are the doubly anharmonic oscillators (DAHO)
$H_{(m, n)} \psi=E_{(m, n)} \psi$,
$H_{(m, n)}=p^{2}+r^{2}+h_{2} p^{4}+g_{2} r^{4}+\ldots+h_{m} p^{2 m}+g_{n} r^{2 n}, \quad h_{m}>0, g_{n}>0$,
investigated by Znojil (1981) and admitting various physical interpretations of the dispersive or 'velocity confining' anharmonicities $p^{4} \ldots$ related originally to the semirelativistic kinetic-energy operator $T(p)=\left(1+p^{2}\right)^{1 / 2}-1+\mathrm{O}\left(p^{2 m+2}\right)$.

When compared with some standard versions of perturbation theory, the nonnumerical mathematical background of the MCF technique has been well developed for $t=1$ only. Nevertheless, the general $t>1 \mathrm{MCF}$ and chain models need not always be treated and interpreted in a semi-numerical context only. Our present intention is to show that some of the analytic $t=1$ techniques may be generalised directly to any $t>1$.

The contents of our paper are organised as follows. In § 2 we review briefly the recurrent MCF idea. Its factorised reinterpretation (Znojil 1983) is recalled and applied to the Feshbach (1958) construction of the effective Hamiltonian in § 3.1. In § 3.2 this is illustrated on the simplest DAHO examples.

In the asymptotic 'quasiconstant' approximation, the elimination of recurrences and the general algebraic construction of $\psi$ and underlying MCF quantities is obtained as our main result in $\S \S 4.1$ and 4.2 . Finally, in $\S 4.3$, the symmetric DAHO example enables us to obtain closed analytic results including, e.g., the explicit $n \gg 1$ asymptotic behaviour of the wavefunctions.

We may summarise that the present method:
(1) clarifies algebraically the factorisable character of the multiterm recurrences and eigenvalue problem in the infinite-dimensional chain models,
(2) simplifies the rigorous analysis of the normalisability of $\psi$ and also, in connection with the semi-numerical MCF technique, the proofs of the MCF convergence,
(3) provides a practical prescription how to construct the optimal and reliable 'smoothed' approximants to the band Hamiltonians in the purely numerical 'cut-off' context, and
(4) enables us to construct various algebraically invertible infinite auxiliary quasiconstant matrices $H_{0}$. As far as simplicity of the unperturbed basis is often preferred even to the smallness of the perturbations $H-H_{0}$, this might help us to achieve suitable rearrangements of perturbation theory (cf e.g. Znojil 1980) and/or to understand its divergences in the future.

## 2. The Feshbach and McF methods

In the general infinite-dimensional linear equation of the non-homogeneous or Schrödinger type,

$$
\begin{equation*}
H \psi=E \psi+\chi \tag{2}
\end{equation*}
$$

we shall consider a non-zero or zero input function $\chi$ such that $\chi=P \chi$, where $P$ denotes a finite-dimensional projector on a 'model' subspace of the whole Hilbert space. It is useful to partition (2) and write

$$
P(H-E) P \psi+P H Q \psi=P \chi, \quad Q=1-P, \quad Q H P \psi+Q(H-E) Q \psi=0
$$

Formally we may eliminate $Q \psi\left[=(E-Q H Q)^{-1} Q H P \psi\right]$ from the second row and obtain the first one in a simple form

$$
\begin{equation*}
\mathscr{H}(E) \varphi=\chi . \tag{3}
\end{equation*}
$$

With $\varphi=P \psi$ and with the definition

$$
\begin{equation*}
\mathscr{H}(E)=P(H-E) P-P H Q[Q(H-E) Q]^{-1} Q H P \tag{4}
\end{equation*}
$$

of the finite-dimensional 'effective' matrix Hamiltonian, this is the essence of the projection method of Feshbach (1958).

Obviously, equation (3) may be solved numerically with arbitrary precision. As a finite-dimensional inversion or eigenvalue problem (with $\chi \neq 0$ or $\chi=0$, respectively) it disguises all the infinite-dimensional algebra of (2) in the definition (4).

In the following text, we shall assume that the operator $H$ has a block-tridiagonal structure in the chosen basis,
$H-E I=\left(\begin{array}{ccc}A_{0} & B_{0} & \\ C_{1} & A_{1} & B_{1} \\ & C_{2} & A_{2} \\ & & \ldots .\end{array}\right) \boldsymbol{B}_{2}, \quad \operatorname{dim} A_{k}=t, \quad k=0,1, \ldots$
As a consequence, definition (4) leads to the partitioned form of the effective Hamiltonians

$$
\mathscr{H}(E)=\mathscr{H}_{(k)}(E)=\left(\begin{array}{cccc}
A_{0} & B_{0} & &  \tag{6}\\
& \ldots & & \\
& C_{k-1} & A_{k-1} & B_{k-1} \\
& & C_{k} & G_{k}
\end{array}\right)
$$

where $P=|0\rangle\langle 0|+|1\rangle\langle 1|+\ldots+|(k+1) t-1\rangle\langle(k+1) t-1|$ and

$$
\begin{equation*}
G_{k}=A_{k}-B_{k} G_{k+1}^{-1} C_{k+1}, \quad k=0,1, \ldots \tag{7}
\end{equation*}
$$

In agreement with Znojil (1980), the MCF recurrences (7) may be initialised by the truncation of $H$, namely, by $G_{N}=A_{N}$ in the limit $N \rightarrow \infty$. This completes the reduction of (2) to its finite-dimensional equivalent (3).

## 3. Chain models and exact effective Hamiltonians

In agreement with $\S 2$, the correct MCF input (exact effective ( $t \times t$ )-dimensional submatrix $G_{k}$ ) reduces (2) to the trivial equation (3) and vice versa, any linear problem (2) with the band matrix Hamiltonian $H$ leads in essence to the infinite-dimensional inversion of $Q(H-E) Q$. Thus, we may replace recurrences (7) by the rigorous non-MCF definition

$$
\begin{equation*}
G_{k}=P_{k} \frac{1}{P_{k} \frac{1}{R_{k}(H-E) R_{k}} P_{k}} P_{k} \tag{8}
\end{equation*}
$$

written in terms of the projectors

$$
\begin{aligned}
& P_{k}=|k t\rangle\langle k t|+|k t+1\rangle\langle k t+1|+\ldots+|(k+1) t-1\rangle\langle(k+1) t-1|, \\
& R_{k+1}=I-P_{0}-P_{1}-\ldots-P_{k}, \quad k \geqslant 0, R_{0}=I .
\end{aligned}
$$

### 3.1. Complete factorisation

Assuming for simplicity that the matrices $B_{k}$ and $C_{k+1}^{\mathrm{T}}$ in (5) are lower triangular, we may evaluate directly the projected inversion $P_{k}\left[R_{k}(H-E) R_{k}\right]^{-1} P_{k}$ in (8) by means of the complete factorisation method (Znojil 1983). In essence, this method is based
on the decomposition of $H-E I$ into the product

$$
\begin{array}{r}
H-E I=\left(\begin{array}{cc}
\rho_{1} & \\
& \\
& \rho_{2} \\
& \\
& \\
K_{\mathrm{U}} & =\prod_{i=1}^{t}\left(\begin{array}{ccc}
1 & \alpha_{1}^{(i)} & 0 \\
0 & 1 & \alpha_{2}^{(i)} \\
0 & 0 & \ldots \\
0 & 0 & 1
\end{array}\right. \\
\alpha_{3}^{(i)}
\end{array}\right) \times K_{\mathrm{U}} \times\left(\begin{array}{cc}
\gamma_{1} & \\
& \gamma_{2} \\
& \\
& \ldots
\end{array}\right) \times K_{\mathrm{L}} \times\left(\begin{array}{lll}
\sigma_{1} & & \\
& \sigma_{2} & \\
& & .
\end{array}\right),  \tag{9}\\
\end{array}
$$

of the $2 t$ two-diagonal and 3-diagonal auxiliary matrices. In our original paper, the detailed algebraic recurrent definitions of the quantities $\rho_{n}, \alpha_{n}^{(i)}, \gamma_{n}, \beta_{n}^{(j)}$ and $\sigma_{n}$ $(i, j=1,2, \ldots, t$, and $n=1,2, \ldots)$ were studied, while here we re-emphasise only their standard initialisation (with $(H)_{M+i, M+j}=0, i, j=1,2, \ldots$ in the 'cut-off' limit $M \rightarrow \infty$ ) and prove the following.

Theorem 1. For the factorised Hamiltonians (9), we may write formally their explicit Feshbach $(N \times N)$-dimensional equivalents in the form (6) with $N=(k+1) t$ and with the effective ( $t \times t$ )-dimensional submatrices

$$
\begin{align*}
& G_{k}=\left(\begin{array}{ccc}
\rho_{N-t+1} & & \\
\cdot & \ddots & \\
& & \rho_{N-1} \\
\rho_{N}
\end{array}\right) \times \prod_{i=1}^{t}\left(\begin{array}{llll}
1 & \alpha_{N-t+1}^{(i)} & & \\
& \cdots & & \\
& 1 & \alpha_{N-1} \\
& & & \\
& & &
\end{array}\right) \times\left(\begin{array}{rll}
\gamma_{N-t+1} & & \\
& \ddots & \\
& & \gamma_{N-1} \\
& & \\
& & \\
& & \\
&
\end{array}\right) \\
& \times \prod_{j=1}^{t}\left(\begin{array}{ccc}
1 & & \\
\beta_{N-t+1}^{(j)} & 1 \\
\cdots & & \\
& \beta_{N-1}^{(j)} & 1
\end{array}\right) \times\left(\begin{array}{ccc}
\sigma_{N-t+1} & & \\
\sigma_{N-t+2} & \\
& \ddots & \\
& & \sigma_{N}
\end{array}\right) . \tag{10}
\end{align*}
$$

Proof. The product form of (9) enables us to write any projection of $H-E$ as the same projected product of the two-diagonal factors. Thus we get (10) after the double use of the formula

$$
\left(\begin{array}{ccc}
1 & \delta_{1} &  \tag{11}\\
& 1 & \delta_{2} \\
& \cdots &
\end{array}\right)^{-1}=\left(\begin{array}{cccc}
1 & -\delta_{1} & \delta_{1} \delta_{2} & \cdots \\
& 1 & -\delta_{2} & \cdots \\
& & \cdots &
\end{array}\right)
$$

in (8).

### 3.2. Example: asymmetric DAHO

In the oscillator basis, the phenomenological daho Hamiltonians $H=H_{(m, n)}$ in (1) have a symmetric and real band-matrix form with $2 t+1$ diagonals, $t=\max (m, n)$. From the formal point of view it is convenient to consider a large index $k$ in (6).

Then we may replace the complicated DAнO matrix elements by their asymptotic expansions. In the leading order, we get the formula

$$
\begin{array}{rlrl}
Q(H-E) Q & =g_{n} N^{n} K(1+\mathrm{O}(1 / N)), & & m<n, \\
& =h_{m} N^{m} J K J(1+\mathrm{O}(1 / N)), & m>n  \tag{12}\\
N=(k+1) t, & Q=R_{k+1}, \quad J=\left(\begin{array}{lll}
-1 & & \\
& 1 & \\
& & -1 .
\end{array}\right),
\end{array}
$$

where $K$ is a constant matrix such that

$$
K=\prod_{i=1}^{t} I_{(+)} \times \prod_{i=1}^{t} I_{(+)}^{\mathrm{T}}, \quad I_{(+)}=\left(\begin{array}{cccc}
1 & 1 & &  \tag{13}\\
& 1 & 1 & \\
& & 1 & 1 \\
& & & \ldots
\end{array}\right)
$$

Hence, we may apply theorem 1 and obtain immediately the effective matrix Hamiltonian $\mathscr{H}(E)$ and equation (3) in the leading-order approximation.

Below we shall use also the symmetric ( $m=n$ ) DAнO chain model as another illustrative example. Due to the exact elimination of its outer diagonal in some bases, its factorisation will be much less trivial.

## 4. Asymptotically smooth chain models

In the following, we shall assume that $Q(H-E) Q$ may be approximated by a constant matrix $K$ (cf example (12)) or by a product $\rho K \sigma$ where $\rho$ and $\sigma$ are some appropriate diagonal matrices introduced in (9), reflecting, e.g., the overall increase of matrix elements with increasing indices, and containing possibly also the error estimates of the type $1+\mathrm{O}(1 / N), N=\operatorname{dim} P, P=1-Q$. Formally, the presence of $\rho$ and $\sigma$ is irrelevant-we may include $\sigma$ in the 'renormalised' vector, $\psi \rightarrow \dot{\sigma} \psi$, and simply multiply equation (1) by the non-singular factor $\rho^{-1}$ from the left.

### 4.1. Algebraic factorisation of $K$

In the real and symmetric case, the factorisation (9) of the constant matrix $K$ may be given the form

$$
\begin{align*}
& K=\left(\begin{array}{lllllll}
a_{t}(t), & a_{t-1}(t), & \ldots & a_{0}(t), & 0, & \ldots & \\
a_{t-1}(t), & a_{t}(t), & a_{t-1}(t), & \ldots & a_{0}(t), & 0, & \ldots \\
\ldots & & & & & &
\end{array}\right)  \tag{14}\\
& =\gamma \times K_{\mathrm{U}} \times K_{\mathrm{L}} \text {, } \\
& K_{\mathrm{U}}=\left(\begin{array}{cccc}
1 & \alpha^{(1)} & 0 & \ldots \\
0 & 1 & \alpha^{(1)} & 0 \ldots \\
& & \ldots &
\end{array}\right) \times\left(\begin{array}{ccccc}
1 & \alpha^{(2)} & 0 & \ldots & \\
0 & 1 & \alpha^{(2)} & 0 & \ldots \\
& & \ldots & &
\end{array}\right) \\
& \ldots \times\left(\begin{array}{ccccc}
1 & \alpha^{(t)} & 0 & \ldots & \\
0 & 1 & \alpha^{(t)} & 0 & \ldots
\end{array}\right)=K_{\mathrm{L}}^{\mathrm{T}},
\end{align*}
$$

with the constant diagonal matrix, $\gamma=1 /\left(\alpha^{(1)} \alpha^{(2)} \ldots \alpha^{(t)}\right)$; with the normalisation of matrix elements $a_{0}(t)=1$, and with the mutually commutative two-diagonal factors in $K_{\mathrm{U}}$. Our intention is to illustrate the general factorisation method of $\S 3$ by its application to the specific matrices (14), i.e., in essence, to define the matrix elements $\alpha^{(t)}, \alpha^{(2)}, \ldots, \alpha^{(t)}$ as functions of $a_{1}(t), a_{2}(t), \ldots, a_{1}(t)$ in an explicit way.

In the simplest $t=1$ case, we have $a_{1}(1)=\alpha^{(1)}+1 / \alpha^{(1)}$ with an elementary inversion. The first non-trivial $t=2$ example

$$
\begin{aligned}
& a_{1}(2)=\alpha^{(1)}+1 / \alpha^{(1)}+\alpha^{(2)}+1 / \alpha^{(2)} \\
& a_{2}(2)=2+\left(\alpha^{(1)}+1 / \alpha^{(1)}\right)\left(\alpha^{(2)}+1 / \alpha^{(2)}\right),
\end{aligned}
$$

inspires us to eliminate the irrelevant $\alpha \rightarrow 1 / \alpha$ ambiguity and to denote $\alpha^{(i)}=\exp \left( \pm \beta_{i}\right)$ or, with $2 \cosh \beta_{i}=X_{(i)}$,

$$
\begin{equation*}
\alpha^{(i)}=\alpha_{( \pm)}^{(i)}=\frac{1}{2} X_{(i)} \pm\left(\frac{1}{4} X_{(i)}^{2}-1\right)^{1 / 2} \tag{15}
\end{equation*}
$$

In general, we then have the following theorem.
Theorem 2. Factorisation (9) of the constant matrix $K$ may be defined by formulae (14) and (15) where the factors $X_{(i)}$ are roots of the polynomial equation

$$
\begin{equation*}
\sum_{k=0}^{i}(-1)^{k} T_{k}(a) X^{t-k}=0 \tag{16}
\end{equation*}
$$

and

$$
T_{k}(a)=\sum_{j=0}^{[k / 2]}(-1)^{j}\binom{t-k+j}{j}\left(1+\frac{j}{t-k+j}\right) a_{k-2 j}(t)
$$

are linear functions of the matrix elements of $K$.
Proof. From the recurrences (14), i.e.

$$
\begin{aligned}
& a_{0}(t)=1 \\
& a_{1}(t)=2 \cosh \beta_{1}+2 \cosh \beta_{2}+\ldots+2 \cosh \beta_{t} \\
& \cdots \\
& a_{k}(t)=a_{k}(t-1)+2 \cosh \beta_{t} a_{k-1}(t-1)+a_{k-2}(t-1), \\
& \cdots \\
& a_{t}(t)=2 \cosh \beta_{t} a_{t-1}(t-1)+2 a_{t-2}(t-1),
\end{aligned}
$$

we infer a possibility to use the ansatz

$$
\begin{equation*}
a_{k}(t)=\sum_{i=0}^{[k / 2]} A_{j}^{(k)}(t) \sum_{\left\{i_{i}\right\}} \cosh \beta_{i_{1}} \cosh \beta_{i_{2}} \ldots \cosh \beta_{i_{k-2 i}} \tag{17}
\end{equation*}
$$

for any $t \geqslant 1$. The unknown coefficients $A_{j}^{(k)}(t)$ may be obtained from the recurrences (14) by insertion. It is documented in table 1 and illustrated by the further $t=3$ example

$$
\begin{aligned}
& a_{0}(3)=1 \\
& a_{1}(3)=2 \cosh \beta_{1}+2 \cosh \beta_{2}+2 \cosh \beta_{3} \\
& a_{2}(3)=3+4\left(\cosh \beta_{1} \cosh \beta_{2}+\cosh \beta_{1} \cosh \beta_{3}+\cosh \beta_{2} \cosh \beta_{3}\right), \\
& a_{3}(3)=2 a_{1}(3)+8 \cosh \beta_{1} \cosh \beta_{2} \cosh \beta_{3} .
\end{aligned}
$$

Now our task is to define the inverse mapping $a_{j} \rightarrow \beta_{i}, 1 \leqslant i, j \leqslant t$.

Table 1. Sample of coefficients $\boldsymbol{A}_{j}^{(k)}(t)$.

|  | 0 | 1 | 2 | 3 |
| :--- | :---: | :--- | :--- | :--- |
| 0 | 1 |  |  |  |
| 1 | 2 |  |  |  |
| 2 | 4 | $t$ |  |  |
| 3 | 8 | $2(t-1)$ | $\binom{t}{2}$ |  |
| 4 | 16 | $4(t-2)$ | $2\binom{t-1}{2}$ |  |
| 5 | 32 | $8(t-3)$ | $4\binom{t-2}{2}$ | $\binom{t}{3}$ |
| 6 | 64 | $16(t-4)$ | $2\binom{t-1}{3}$ |  |
| 7 | 128 | $32(t-5)$ |  |  |

Our first result concerns the compact and explicit formula

$$
\begin{equation*}
A_{j}^{(k)}(t)=2^{k-2 i}\binom{t-k+2 j}{j} \tag{18}
\end{equation*}
$$

which may be verified by mathematical induction and a certain amount of patient elementary manipulations. As a consequence, the mapping (17) may be reinterpreted as a composition of nonlinear transformations $\cosh \beta_{k} \rightarrow S_{i}\left(\cosh \beta_{1}, \ldots, \cosh \beta_{t}\right)$ with a pair of two separate linear mappings
$a_{2 i}(t)=\sum_{k=0}^{i}\binom{t-2 k}{i-k} 4^{k} S_{2 k}\left(\cosh \beta_{1}, \cosh \beta_{2}, \ldots, \cosh \beta_{t}\right), \quad i=1,2, \ldots,[t / 2]$,
and

$$
\begin{gather*}
a_{2 i-1}(t)=\sum_{k=1}^{i}\binom{t-2 k+1}{i-k} 2^{2 k-1} S_{2 k-1}\left(\cosh \beta_{1}, \cosh \beta_{2}, \ldots, \cosh \beta_{t}\right)  \tag{20}\\
i=1,2, \ldots,[(t+1) / 2]
\end{gather*}
$$

where

$$
\begin{gather*}
S_{0}\left(x_{1}, x_{2}, \ldots, x_{n}\right)=1, \quad S_{1}\left(x_{1}, x_{2}, \ldots, x_{n}\right)=x_{1}+x_{2}+\ldots+x_{n}, \ldots, \\
S_{n}\left(x_{1}, x_{2}, \ldots, x_{n}\right)=x_{1} x_{2}, \ldots, x_{n}, \tag{21}
\end{gather*}
$$

are the ordinary symmetric functions defined by the identity

$$
\begin{equation*}
\sum_{m=0}^{N} x^{N-m} S_{m}\left(x_{1}, x_{2}, \ldots, x_{N}\right)=\prod_{i=1}^{N}\left(x+x_{i}\right) . \tag{22}
\end{equation*}
$$

The matrices of the linear systems (19) and (20) are triangular and easily invertible (appendix). Hence, we may also invert the mapping $a_{i}(t) \leftarrow S_{i}$,

$$
\begin{equation*}
S_{m}\left(\cosh \beta_{1}, \ldots, \cosh \beta_{t}\right)=2^{-m} T_{m}(a), \quad m=0,1, \ldots, t \tag{23}
\end{equation*}
$$

Finally, the definition (22) of the symmetric functions $S_{m}$ implies that the polynomial equation

$$
\sum_{m=0}^{t}(-1)^{m} 2^{-m} T_{m}(a)(\cosh \beta)^{t-m}=0
$$

determines all the factors $\alpha^{(i)}$ in (14) in a purely algebraic way.

### 4.2. Wavefunctions

Obviously, in terms of the $\beta$ or $X$ variables defined by theorem 2 , the purely algebraic definition of the asymptotic effective Hamiltonian (6) is an easy consequence of theorem 1. Nevertheless, rather surprisingly, the matrix elements of $K^{-1}$ are not unique in general. For example, we get $G_{k}^{(1)}=1 / \alpha_{(+)}^{(1)}$ and $G_{k}^{(2)}=1 / \alpha_{(-)}^{(1)}=\alpha_{(+)}^{(1)}$ for $t=1$.

This ambiguity is not encountered in the asymptotic $m \neq n$ DAHO examples (12) where $\alpha_{( \pm)}^{(i)}=1, i=1,2, \ldots, t$. In the other, non-degenerate chain models, we may prove the following.

Theorem 3. Provided that $\sigma_{k} / \sigma_{k-1}=1+\mathrm{O}(1 / k), k \gg 1$, and $X_{(i)}>2, X_{(i)} \neq X_{(j)}, i, j=1$, $2, \ldots, t$, the norm of the wavefunctions.

$$
\|\psi\|=\left(\sum_{n=0}^{\infty}\langle\psi \mid n\rangle\langle n \mid \psi\rangle\right)^{1 / 2},
$$

is finite if and only if we choose $\alpha^{(i)}=\alpha_{(-)}^{(i)}$ in (15).
Proof. Since $\left|\alpha_{(+)}^{(i)}\right|>1$ and $\left|\alpha_{(-)}^{(i)}\right|<1$, we may neglect the higher-order corrections and write the definitions

$$
\begin{aligned}
& Q K Q \psi=-Q K P \psi=\left(\begin{array}{c}
c_{1} \\
\vdots \\
c_{t}
\end{array}\right), \\
& Q \psi=\frac{1}{Q K Q}\left(\begin{array}{c}
c_{1} \\
\vdots \\
c_{t}
\end{array}\right)=K_{\mathrm{L}}^{-1} K_{\mathrm{U}}^{-1}\left(\begin{array}{c}
c_{1} \\
\vdots \\
c_{t}
\end{array}\right)=K_{\mathrm{L}}^{-1}\left(\begin{array}{c}
d_{1} \\
\vdots \\
d_{t}
\end{array}\right) \\
& \quad=\left(\begin{array}{c}
1 \\
-\alpha^{(1)}-\alpha^{(2)}-\cdots \\
\cdots
\end{array}\right)\left(\begin{array}{c}
(t) \\
\cdots \\
\vdots \\
d_{t}
\end{array}\right) .
\end{aligned}
$$

Here, the explicit form of the matrix elements may be inserted, which gives
$\langle N+m-1 \mid \psi\rangle=\sum_{k=1}^{\min (m, t)}(-1)^{k+m} d_{k} \sum_{i=1}^{t} \frac{\left(\alpha^{(i)}\right)^{m+t-k-1}}{\prod_{j \neq i}\left(\alpha^{(i)}-\alpha^{(j)}\right)}, \quad m=1,2, \ldots$
This is one of our most interesting formulae and implies that, for $m \gg 1$, we have an estimate

$$
\begin{equation*}
|\langle m \mid \psi\rangle /\langle m-1 \mid \psi\rangle| \doteq \max _{1 \leqslant i \leqslant i}\left|\alpha^{(i)}\right|, \quad m \gg 1 . \tag{25}
\end{equation*}
$$

Hence, the infinite series $\|\psi\|$ converges or diverges as a geometric series.

To avoid confusion, we must re-emphasise that the above energy-independent construction of the normalisable vector $\psi$ is not contradictory since this vector satisfies the Schrödinger equation at the physical energies only.

### 4.3. Example: symmetric DAHO

After an appropriate scaling of coordinates $r \rightarrow \nu r, p \rightarrow(1 / \nu) p$, an approximate cancellation of diagonals in the daho $m=n$ matrix $H_{(m, n) i j}$ may be achieved in the $i, j \gg 1$ asymptotic region. Thus, instead of (12), we get the leading-order formula

$$
\begin{array}{rlr}
H_{(m, n) i j}-E \delta_{i j}=\rho_{i} \tilde{K}_{i j} \sigma_{j}, & \quad i, j \gg 1,  \tag{26}\\
=1, & i=j, \quad \tilde{K}_{i j}=\binom{2 m}{m+i-j} \pi_{i j}, \\
\delta_{i j}=0, & i \neq j, \\
=1, & & |i-j|=\text { even, } \\
\pi_{i j} & =0, & \\
\pi_{i}, j \mid=\text { odd }, &
\end{array}
$$

where $\rho_{i}$ and $\sigma_{i}$ contain again the irrelevant $i$ and $j$ dependence of $H_{i j}$ and possibly also the explicit non-diagonal error estimates $1+\mathrm{O}(1 /(i+j))$.

The factor $\pi_{i j}$ is a projector on the even- or odd-indexed subspaces so that it is sufficient to consider just one of the two orthogonal and identical infinite matrices $K$ with the smaller parameter $t=t_{\mathrm{as}}=$ integer part of $\left(\frac{1}{2} m\right)$. Thus with the matrix elements

$$
\begin{equation*}
K_{i j}=\binom{2 m}{m+2 i-2 j}, \quad i, j=1,2, \ldots \tag{27}
\end{equation*}
$$

we may use theorems 1 and 2 .
One of the most obvious ways to eliminate the $\alpha \rightarrow 1 / \alpha$ ambiguity without any recourse to theorem 3 is to return to the original MCF interpretation of $G_{k}$. We do not get any contradiction-in the simplest quartic-quartic example with $m=n=$ $2 t_{\mathrm{as}}=2$, the continued-fractional result

$$
\begin{equation*}
G_{k} \sim 6-\frac{1}{6-\frac{1}{6-\ldots}}=3+2 \sqrt{2} \doteq 5.828 \tag{28}
\end{equation*}
$$

is easily summable and eliminates the plus-sign choice in (15). Similarly, for $m=n=3$, we get $t_{\text {as }}=1$ and

$$
\begin{equation*}
G_{k} \sim 20 / 6-\frac{1}{20 / 6-\frac{1}{20 / 6-\ldots}}=3 \tag{29}
\end{equation*}
$$

Starting from $m=n=4$, the corresponding algebra may be combined with the numerical tests as well.

In the light of theorem 3, criterion (25) specifies now the upper $P \psi$-independent estimate of the rate of convergence of $\|\psi\|$. Up to the sufficiently high-order potentials $V=H_{(m, n)}-p^{2} \sim p^{18}+\ldots+r^{18}$ with $t_{\mathrm{as}} \leqslant 4$, this is still a non-numerical task in prin-ciple-we may obtain $\alpha^{(i)}$ by algebraic means. For convenience, the corresponding $t \leqslant 4$ linear functions $T_{k}(a)$ to be used in (16) are listed here in table 2.

In table 3 the explicit exact values of $\cosh \beta_{i}$ are given up to the duodecadicduodecadic DAHO case. We may notice that the assumptions of theorem 3 appear to be satisfied so that, in contrast to the asymmetric DAHO examples, the corresponding

Table 2. First five functions $T_{i}(a)=2 S_{i}\left(\cosh \beta_{1}, \ldots, \cosh \beta_{t}\right)$.

| $T_{0}(a)=1$ | $t \geqslant 0$ |
| :--- | :--- |
| $T_{1}(a)=a_{1}$ | $t \geqslant 1$ |
| $T_{2}(a)=a_{2}-t$ | $t \geqslant 2$ |
| $T_{3}(a)=a_{3}-(t-1) a_{1}$ | $t \geqslant 3$ |
| $T_{4}(a)=a_{4}-(t-2) a_{2}+t(t-3) / 2$ | $t \geqslant 4$ |

Table 3. Leading-order asymptotic factorisation of the symmetric DAHO operator $H=$ $p^{2 m}+\ldots+r^{2 m}$.

| $m$ | $i$ | $\frac{1}{2} X_{(1)}$ | $\alpha^{(1)}$ |
| :--- | :--- | :--- | :--- |
| 2 | 1 | 3 | $0.171572875 \ldots$ |
| 3 | 1 | $5 / 3$ | $0.333333333 \ldots$ |
| 4 | 1 | $7+4 \sqrt{2}$ | $0.039566129 \ldots$ |
| 5 | 2 | $7-4 \sqrt{2}$ | $0.446462692 \ldots$ |
|  | 1 | $3+4 / \sqrt{5}$ | $0.105572808 \ldots$ |
| 6 | 2 | $3-4 / \sqrt{5}$ | $0.527864045 \ldots$ |
|  | 1 | $15+8 \sqrt{3}$ | $0.017332380 \ldots$ |
|  | 2 | 3 | $0.171572875 \ldots$ |
|  | 3 | $15-8 \sqrt{3}$ | $0.588790706 \ldots$ |

wavefunctions lie mostly within a small model space. Also the norm $\|\psi\|$ converges extremely quickly. All these features of the symmetric DAHO system are compatible with the preceding numerical results (Znojil 1981).

## Appendix. Inversion of the triangular matrices

For the triangular matrices

$$
\left(\begin{array}{cccc}
1 & & &  \tag{A1}\\
a_{1} & 1 & & \\
b_{1} & a_{2} & 1 & \\
c_{1} & b_{2} & a_{3} & 1 \\
\ldots & & &
\end{array}\right)
$$

the inverse

$$
\left(\begin{array}{ccc}
1 & &  \tag{A2}\\
\alpha_{1} & 1 & \\
\beta_{1} & \alpha_{2} & 1 \\
\ldots & 1
\end{array}\right)
$$

is defined by the obvious relations

$$
\begin{align*}
& a_{i}+\alpha_{i}=0, \quad b_{i}+a_{i+1} \alpha_{i}+\beta_{i}=0, \\
& c_{i}+b_{i+1} \alpha_{i}+a_{i+2} \beta_{i}+\gamma_{i}=0, \quad \ldots, \quad i=1,2, \ldots, \tag{A3}
\end{align*}
$$

and it may be given the explicit Kramer form (Korn and Korn 1968)

$$
\begin{align*}
\alpha_{i} & =-\operatorname{det}\left(a_{i}\right),  \tag{A4}\\
\gamma_{i} & =-\operatorname{det}\left(\begin{array}{ccc}
a_{i} & 1 & 0 \\
b_{i} & a_{i+1} & 1 \\
c_{i} & b_{i+1} & a_{i+2}
\end{array}\right), \quad \ldots\left(\begin{array}{cc}
a_{i} & 1 \\
b_{i} & a_{i+1}
\end{array}\right),
\end{align*}
$$

In particular, for the matrix elements (18),
$a_{i}=\binom{T-2 i+2}{1}, \quad b_{i}=\binom{T-2 i+2}{2}, \quad c_{i}=\binom{T-2 i+2}{3}, \quad \ldots$,
we may evaluate (A4)

$$
\begin{align*}
& \alpha_{i}=-(T-2 i+2), \quad \beta_{i}=+\frac{1}{2}(T-2 i+2)(T-2 i-1), \\
& \gamma_{i}=-\frac{1}{6}(T-2 i+2)(T-2 i-2)(T-2 i-3), \quad \cdots, \tag{A6}
\end{align*}
$$

and, changing the notation to

$$
\alpha_{i}=B_{i}^{(1)}(T), \quad \beta_{i}=B_{i}^{(2)}(T), \quad \gamma_{i}=B_{i}^{(3)}(T), \ldots
$$

we arrive at the final formula

$$
\begin{equation*}
B_{i}^{(k)}(T)=(-1)^{k} \frac{T-2 i+2}{T-2 i+2-k}\binom{T-2 i-k+2}{k}, \quad i=1,2, \ldots \tag{A7}
\end{equation*}
$$

which defines the integer coefficients needed in (19) and (20) in §4.1.

## References

Feshbach H 1958 Ann. Phys., NY 5357
Graffi S and Grecchi V 1975 Lett. Nuovo Cimento 12425
Halliday I G and Suranyi P 1980 Phys. Rev. D 211529
Haydock R 1980 Comp. Phys. Commun. 2011
Korn G A and Korn T M 1968 Mathematical Handbook (New York: McGraw-Hill)
Simon B 1969 Ann. Phys., NY 5876
Znojil M 1980 J. Math. Phys. 211629

- 1981 Phys. Rev. D 24903
- 1983 J. Phys. A: Math. Gen. 163313

