

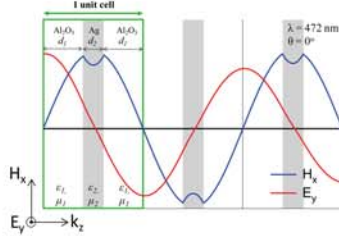
THE SPATIAL DISPERSION EFFECT IN STRATIFIED METAL-DIELECTRIC METAMATERIAL

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I. Introduction

The exotic electromagnetic response in typical metamaterial is attributed to the unusual inhomogeneous electromagnetic distribution in unit cell at particular frequencies. Therefore, the unit cell should small enough compare to wavelength which replaces the role of atom and molecule in conventional materials. But special unit cell design is not the only way to realize an inhomogeneous electromagnetic distribution.



Stratified metal dielectric metamaterial, or SMDM, is an artificial metal dielectric composite that consists of silver (30 nm) sandwiched by identical alumina (60 nm). Usually the optical response in optical frequency range is simply ascribed to the 1-dimensional photonic crystal effect, but we go to one further step to generalize the concept of effective medium to the structure with much larger unit cell size. We regard the artificial structure to be a hypothetical uniform material, which can be described in terms of effective permittivity and effective permeability.

II. Retrieval Procedure

The Maxwell Equation in a unit cell:

$$\oint_C \vec{E} \cdot d\vec{l} = -\frac{\partial}{\partial t} \iint_S \vec{B} \cdot \vec{n} dS$$

$$\oint_C \vec{H} \cdot d\vec{l} = \frac{\partial}{\partial t} \iint_S \vec{D} \cdot \vec{n} dS$$

The Bloch Boundary Condition:

$$\vec{E}_x \left[\left((n+1/2)d \right) \right] = \vec{E}_x \left[d/2 \right] e^{i(n\theta + \theta/2)}$$

$$\vec{H}_y \left[(nd) \right] = \vec{H}_y \left[0 \right] e^{in\theta}$$

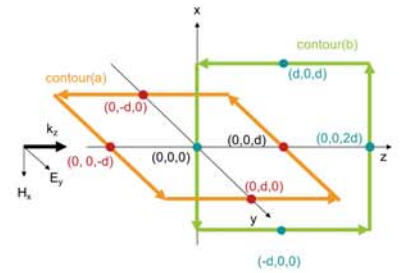
Where θ is the phase advance across one cell

The dispersion equation in SMDM can be described as:

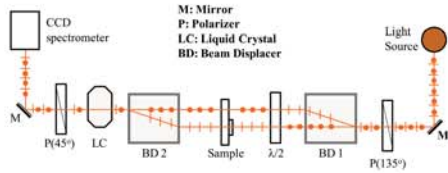
$$\sin(\theta/2) = \omega d \sqrt{\mu \epsilon} / 2$$

The averaging impedance is described as:

$$Z = \sqrt{\frac{\mu}{\epsilon}} \cos(\theta/2)$$

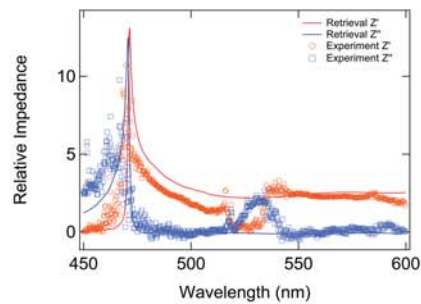


III. Experimental set-up



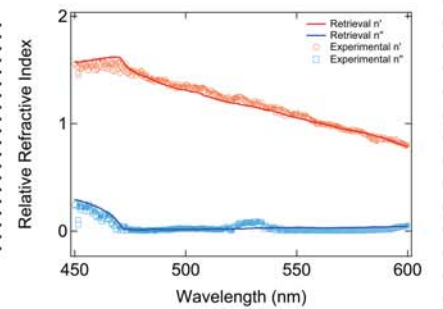
Schematic set up phase measurement. Polarization directions in two arms in the Mach-Zehnder interferometer are indicated by bars (horizontal) and dots (vertical). The second polarizer is set to be orthogonal to the first one. Additional retardation to the reference arm is introduced by liquid crystal variable retarder to recover the cross polarization condition for every wavelength.

IV. Results and Discussions



- The field distribution graph was calculated for a sample with many unit cells in order to suppress the reflected wave from the interface.
- The cusp-like features for H-field at the dielectric-metal interface is explained by Ampere-Maxwell equation with different sign of permittivity in Ag and alumina.
- The largest permeability is 17 for the experiment and 20 for the numerical calculation. The origin of permeability is due to inhomogeneity of the field distribution in the unit cell. The condition occurs when the half wavelength of the light in the structure matches to the unit cell size, when the microscopic magnetic field is concentrated in the metal. Interpretation of effective medium is valid as far as the incident wavelength is so long that no diffraction is allowed.
- The permittivity and permeability resonance around 520 nm in experimental data is ascribed to the Fabry-Perot interference. Such structure does not show in calculation data unless we introduce inequivalent unit cell.

□ When we apply Bloch boundary system in our study, the three cases are observed in different wave modes in SMDM. From 460 to 470 nm, the real refractive index is constant. It means the phase advance reach $\pi/2$ and the wave is resonant crystal band-gap regime. Above 470 nm, all modes correspond to propagating modes. The evanescent modes cannot be observed in the graph, it occurs in higher wavelength regime. In our case, we limits the observation just in optical regime.



V. Conclusion

We have investigated an optical response of metal dielectric multilayers with sub-wavelength period. Our calculation predicted enhancement of permeability, while experimental data of three periods SMDM confirmed it. Imperfections in fabrication process responsible to create another resonance correspond with Fabry-Perot interference. There are no ordinary materials exhibiting magnetic response in optical frequency. By combining each layer as SMDM, which has effective μ_{eff} , we succeeded to create high fascinate magnetic response in this structure. SMDM is simple but useful and pedagogic example to understand the origin of magnetic response in artificial structure consisting of non-magnetic materials.

References

- [1] R. Liu et. al., PRE 76 026606 (2007).
- [2] R. Watanabe et. al., Phys. Status Solidi B 245 2696 (2008).
- [3] P. B Johnson et al. PRB 6 4370 (1972)
- [4] D. R. Smith et al., PRB 65 195104 (2002)
- [5] A. S. Vioktalamo et al., e-print PNFA: 10.1016/j.photonics.2011.08.005

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1. INTRODUCTION

Primary Process of Photosynthesis

Light Harvesting (LH) Function

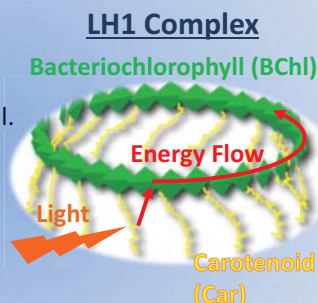
Efficient use of solar light energy:
Ultrafast energy transfer from Carotenoid to Bacteriochlorophyll.

Photoprotective Function

Preventing damage to the cell:
Carotenoid absorbs strong light energy.

Main theme:

Dynamics of energy transfer in Light Harvesting.



2. OBJECTIVE

Hot S_1 is Relaxed from S_2 and Involved with Energy Transfer in LH.

More efficient energy transfer than S_1 . Why does it happen?

What's the hot S_1 state?

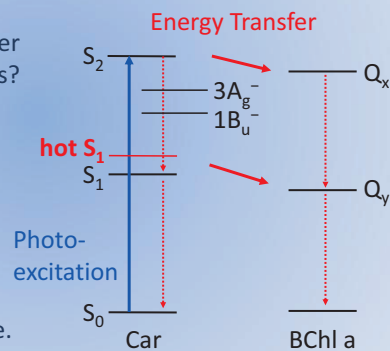
Vibrational state of S_1 ?

Electronic state?

($3A_g^-$ or $1B_u^-$)

S_1 is optical forbidden.

It is not very easy to examine.

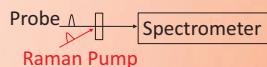


3. METHODS

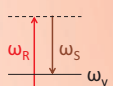
Femtosecond Stimulated Raman Spectroscopy (FSRS)

Stimulated Raman Spectroscopy

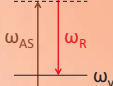
Detection of Transmittance Change



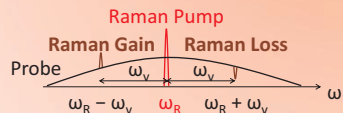
Raman Gain (Stokes)



Raman Loss (Anti-Stokes)

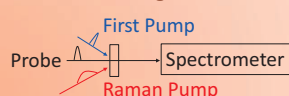


Raman Gain/Loss in Probe Continuum

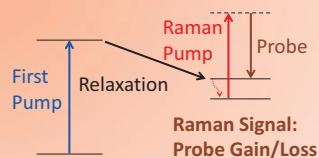


FSRS

Raman Scattering of Transient State



Schematic of FSRS



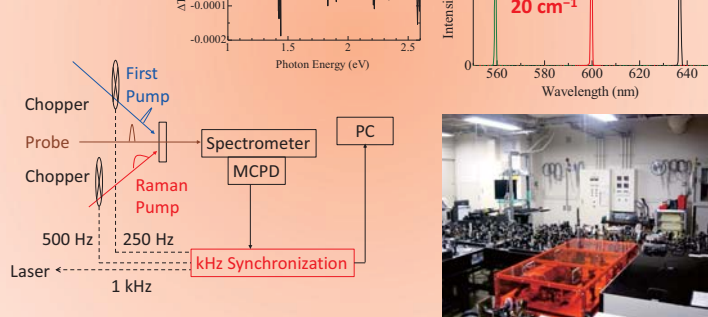
4. PERFORMANCE

FSRS Sensitivity and Spectral Resolution

kHz Synchronized Detection System

Sensitivity:

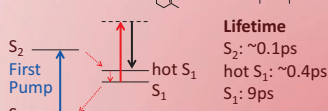
$\Delta T/T: 10^{-5} - 10^{-4}$



5. RESULTS

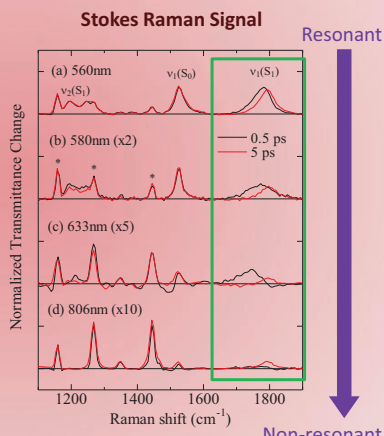
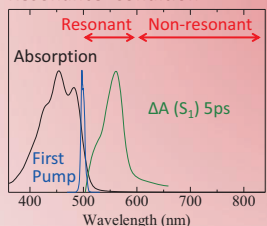
FSRS of β -Carotene Resonant to S_1

β -Carotene



Lifetime
 S_2 : ~ 0.1 ps
hot S_1 : ~ 0.4 ps
 S_1 : 9 ps

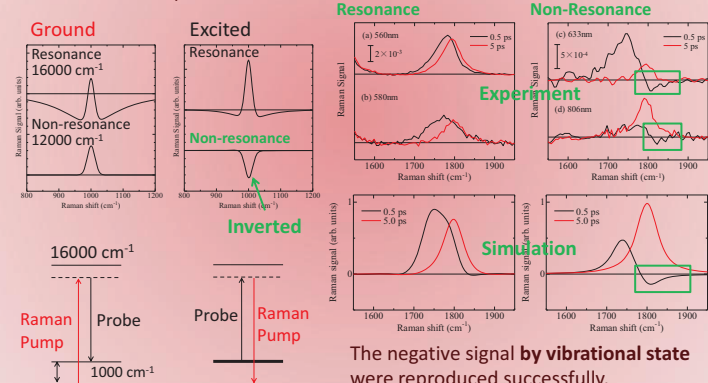
Absorption Spectra and Resonance Condition



6. DISCUSSION

Simulating Negative Signal of $\nu_1(S_1)$ Mode

Stokes Raman Signal Simulations of Vibrational Ground/Excited State



The negative signal by vibrational state were reproduced successfully.

Hot S_1 has been confirmed to the ν_1 vibrational state of S_1 — not to the other electronic state.

Electronic State and Superconductivity of CeRhSi₃

Department of Physics, Tohoku University

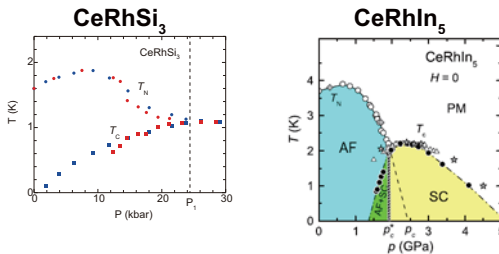
Hiroki Iida

Introduction

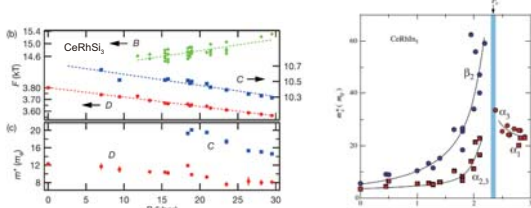
Heavy-fermion superconductivity appears in the vicinity of the border of magnetism in which a magnetic transition temperature goes zero. CeRhIn₅ is one of the typical heavy-fermion superconductors. The superconducting transition temperature T_C becomes maximum at the pressure where the antiferromagnetic order vanishes. An application of pressure suppresses the antiferromagnetic order and simultaneously induces the superconductivity. The effective mass of the conduction electron diverges at this pressure [1]. Interplay between the emergence of superconductivity and the critical phenomenon of the magnetism is an interesting issue for the heavy-fermion system.

CeRhSi₃ is a heavy-fermion superconductor found in 2005 [2]. Although pressure-temperature phase diagram of CeRhSi₃ is similar to that of CeRhIn₅, a critical behavior has not been observed thus far. In order to verify whether or not the effective mass diverges at a critical pressure, we have measured the electrical resistivity under pressure and magnetic field.

Pressure-induced heavy-fermion superconductivity



- Pressure-temperature phase diagram of CeRhSi₃ [2] is similar to that of CeRhIn₅ [3].



- No divergence of effective mass [4]
- Effective mass diverges critical pressure [1]

References

- [1] H. Shishido et al., J. Phys. Soc. Jpn. **74** (2005) 1103.
- [2] N. Kimura et al., J. Phys. Soc. Jpn. **76** (2007) 051010.
- [3] G. Knebel et al., J. Phys. Soc. Jpn. **77** (2008) 114704.
- [4] T. Terashima et al., Phys. Rev. B **76** (2007) 054506.

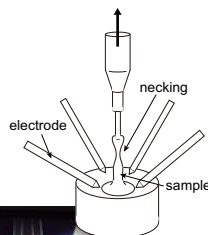
Sample and Measurement

Sample preparation

- The single crystal of CeRhSi₃ was grown by the Czochralski pulling method in a tetra-arc furnace.
- Residual resistivity ratio (RRR) > 100

Application of pressure system

- Clamped piston cylinder cells
- Pressure transmitting medium : equal mixture of n- and i-propanol

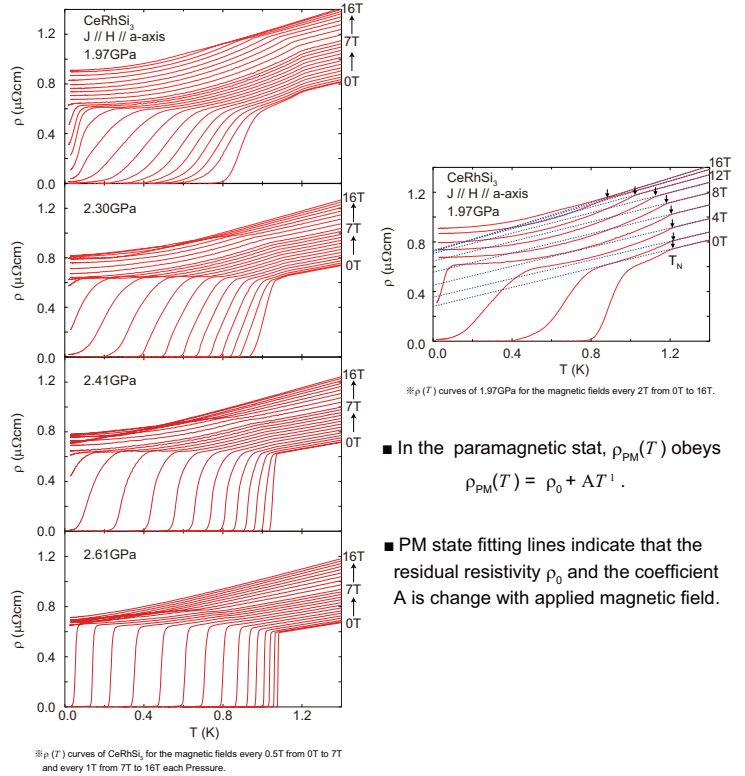


Measurement

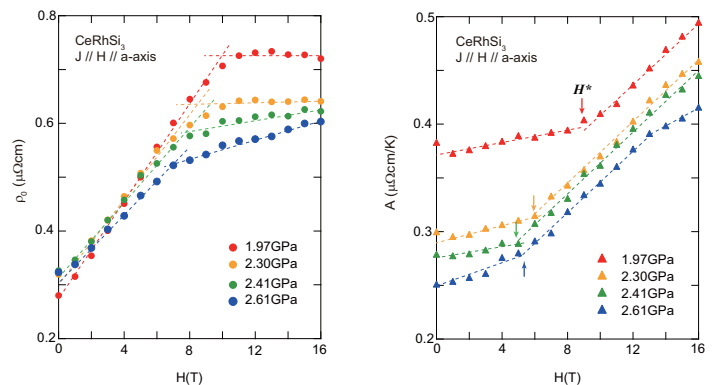
- Low ac four-terminal method



Resistivity under magnetic field for 1.97, 2.30, 2.41 and 2.61 GPa



Magnetic field dependences of residual resistivity ρ_0 and the coefficient A



- The $d\rho_0/dH$ changes at 7~9 T each pressure.
⇒ We suppose that change of the $d\rho_0/dH$ is caused by the valence transition or change of electronic structure.
- The coefficient A exhibit monotonic increase with applied magnetic field.
⇒ No signatures of quantum criticality !!

Summary

- 1) In the paramagnetic stat of CeRhSi₃, $\rho_{PM}(T)$ obeys $\rho_{PM}(T) = \rho_0 + AT^{-1}$.
- 2) The coefficient A does not diverge under magnetic field. Therefore critical behavior of CeRhSi₃ has not been observed.
- 3) No critical behavior of CeRhSi₃ is consistent with dHvA results.

Yoshihisa Ishikawa, Terutoshi Sakakura, Takahiro Nakano, Hiroyuki Kimura, Yukio Noda, Yasuyuki Takenaka¹, Syunji Kishimoto²

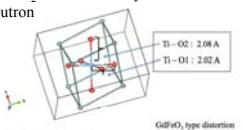
IMRAM, Tohoku Univ., Hokkaido University of Education¹, KEK-Photon Factory²

Introduction

YTiO_3 is well known as one of orbital-ordering system materials. The orbital-ordering phenomenon of YTiO_3 have many studied the both side of theoretical and experimental method. The crystal structure of YTiO_3 belong to a perovskite, Pbnm (Space group No.62 Pnma). The wycoff position of Ti atom is 4a and the site symmetry is -1. Figure shows the crystal structure of yttrium titanate. YTiO_3 undergoes the phase transition from paramagnetic to ferromagnetic at $T_c=28\text{K}$.

In the cubic crystal field, one 3d electron of the Ti^{3+} ion may occupy the triply degenerate t_{2g} states. The Jahn – Teller distortion of TiO_6 octahedron in YTiO_3 splits the energy levels of t_{2g} state. In the t_{2g} energy levels of TiO_6 octahedron, $\text{Ti}3d$ has $|yz\rangle$ and $|zx\rangle$ orbitals. The electron of the orbitals $c_1|zx\rangle - c_2|yz\rangle$ and $c_1|zx\rangle + c_2|yz\rangle$, where coefficients of c_1 and c_2 are satisfied $c_1^2 + c_2^2 = 1$. The study to be determined coefficients c_1 and c_2 is performed by Polarized neutron diffraction (PND) [1], the resonant X-ray scattering (RXS) [2], NMR [3], and the X-ray Magnetic Diffraction (XMD) [4].

The aim of this experiment is to perform the accurate crystal structure analysis by Synchrotron X-ray Diffraction to obtain the coefficients c_1 and c_2 .



[1] J. Akimitsu, et al, JPSJ **70**, 3475 (2001), [2] H. Nakao, et al, PRB **66**, 184419 (2002)
[3] M. Itoh, et al, JPSJ **68**, 2783 (1999), [4] M. Itoh, et al, J. Phys. Chem. Solid **65**, 1993 (2004)

Experiment



X-ray 4-circle Diffractometer @ KEK PF BL-14A (Tsukuba)

- Perpendicular wiggler type x-ray source
- Monochrometer Si (1 1 1) : Wave Length : 0.752Å / 0.850Å
- $2\theta < 122^\circ$ ($Q < 1.163$)
- Detector : Si - APD [Avalanche PhotoDiode] (*)

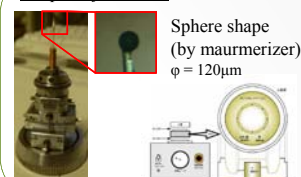
Temperature control

- He-gas flow

Temp	R(F)	meas pt.	used pt.
23K	4.9	1339	1140
32K	3.08	2750	2430
55K	3.2	268	213
85K	2.44	4833	4353
145K	2.52	1683	1919
293K (RT)	1.04	3198	2892

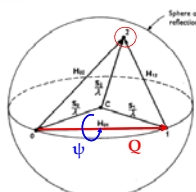
(*) S. Kishimoto, et al, PRL **85**, 1831 (2000)

Sample Information



Result

Avoidance of Multiple Scattering



Simultaneous reflections means several Bragg peaks ride on the Ewald sphere.

To avoid the simultaneous reflections, the setting angle would be rotated around Q-vector which is called "ψ-axis"

In this measurement, the avoidance of multiple scattering performed the program "MDC++" which is coded by Dr. Sakakura.

Fig. 1. Geometric representation of simultaneous reflections in the reciprocal lattice. Points 1, 2 and 3 lie on the sphere of reflection. The unit vectors S_1 , S_2 and S_3 define the direction of the incident, primary and secondary beams.
R. M. Moon, et al, Acta Cryst **17**, 805 (1964)
K. Tanaka, et al, Acta Cryst **A66**, 438 (2010)

Crystal Structure Analysis

Lp Correction : Lorentz Corr. = 2θ
linear polarization $p = 1$.

Absorption Correction : program by DABEX

Extinction Correction. : Becker – Coppens, Type 1, Gaussian Model

$$E_k = \sqrt{1 + 2x + \frac{0.58 + 0.48 \cos 2\theta + 0.24 \cos^2 2\theta}{1 + (0.02 - 0.025 \cos 2\theta)x}}$$

$$x = \frac{7.9406 \times 10^5 \lambda^2 E_0^2 \sin^2 2\theta}{V_0^2 \sin 2\theta (\cos 2\theta_m + \cos^2 2\theta)}$$

P. Becker & P. Coppens, Acta Cryst **A30**, 129 (1974)

Multipole Refinement

$$\rho(\mathbf{r}) = \frac{\rho_c(\mathbf{r})}{r} + \kappa^3 P_v \rho_v^2(\kappa \cdot \mathbf{r}) + \sum_{l=0}^{l_{max}} \kappa^3 R_l(\kappa \cdot \mathbf{r}) \sum_{m=\pm l} P_m Y_{lm}(\theta, \phi)$$

Core Electron valence Electron Spherical Harmonic Function

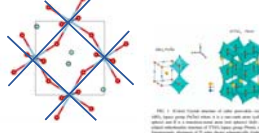
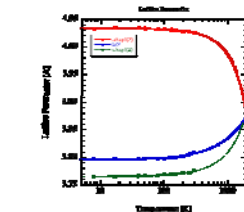
Radial distribution function

$$R_l(r) = \frac{\xi_l^{n_l+3}}{(n_l+2)!} r^{n_l} e^{-\xi \cdot r}$$

Radial distribution function is based on Slater-like Hartree-Fock-Roothaan algebraic equation. (Clementi Table)

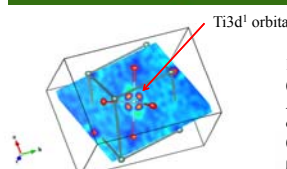
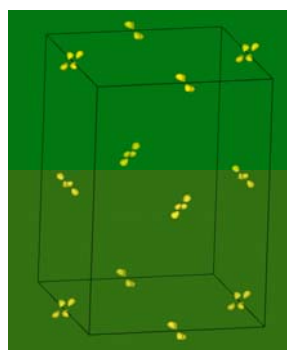
N. K. Hansen & P. Coppens, Acta Cryst **A34**, 909 (1978)
E. Clementi & C. Roetti, ATOMIC DATA AND NUCLEAR DATA TABLES **14**, 177 – 478 (1974)
M. Dusek, et al, Journal of Physics **226**, 012014 (2010)

Framework of Lattice & Octahedron



N. N. Kovaleva, et al, PRB **79**, 045114 (2009)

Electron Density



Deformation Density :

$$\Delta \rho_{model}(\mathbf{r}) = \sum_{i \text{ all atoms}} \{ P_{i,c} \rho_{core}(\mathbf{r}) + P_{i,v} \kappa^3 \rho_{valence}(\kappa_i \mathbf{r}) \}$$

$$+ \sum_{i=0}^{l_{max}} \kappa_i^3 R_{i,l}(\kappa_i r) \sum_{m=0}^l P_{i,lm} d_{lm}(\mathbf{r}/r) \} - \rho_{reference}(\mathbf{r})$$

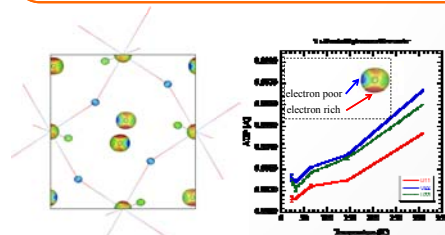
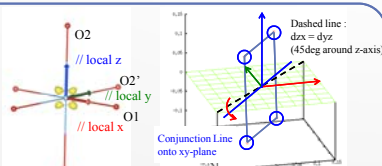


FIGURE.
(Left) Deformation Density of YTiO_3 by multipole refinement at room temperature. And 2-dimensional Differential Fourier Mapping on $\text{Ti}3d$ orbital. Electron density level of $0.8e/\text{\AA}^3$.
(Right) Contour valence electron density mapping on observed Fourier transform density ρ_{obs} . The Blue and Red color means the poor and rich electron number.

Discussion

To determine the coefficients c_1 and c_2 , Ti-O vector has a projection onto a local coordinate system. In this case, Ti-O1, Ti-O2, and Ti-O2 is transformed x, y, and z onto a local orthogonal coordinate. The valence electron density of $\text{Ti}3d$ orbitals can estimate the position by calculation of center of weight. In the crystal field theory, these 3d1 wave functions is representation of linear combination between $|zx\rangle$ and $|yz\rangle$. It means the rotation around local z-axis.

Since Ti valence electron density position is tilted 9 degree than $d_{zx} = d_{yz}$ line (45 degree: In the case of $d_{zx} = d_{yz}$, $c_1 = \sqrt{2}/2 = 0.71$), the coefficients of $\text{Ti}3d$ of this experiment evaluates $c_1=0.81$. This result is consistent with NMR measurement.



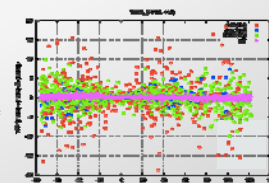
c_1	measurement	Reference
0.69	NMR	M.Itoh, JPSJ 68 , 2783
0.71	RXS, XMD	H.Nakao, PRB 66 , 184419
0.77	PND	J.Akimitsu, JPSJ 70 , 3475
0.8	NMR	T.Kiyama, JPSJ 74 , 1123
0.81	This Experiment	
0.84	Magnetic Compton	F.Iga, PRL 93 , 257207

Summary

In the crystal field theory, $\text{Ti}3d$ of the octahedron in YTiO_3 has the 2 energy levels $|dxz\rangle$ and $|dyz\rangle$. To determine the coefficients of linear combination c_1 and c_2 , the accurate crystal structure analysis performed by Synchrotron X-ray Diffraction.

Estimation for electron density involved $\text{Ti}3d$ t_{2g} orbitals, it succeeds the result for R-Factor within 1% by multipole refinement. And then, the coefficient c_1 was obtained 0.81.

In this measurement, the environment of a low temperature was not optimized. In the next step, we try to improve the measurement methods for a low temperature below T_c .



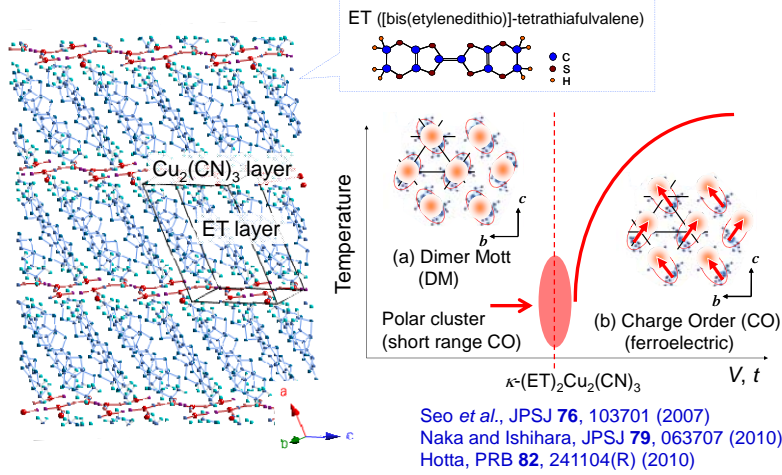
Terahertz Time Domain Spectroscopy of Dimer Mott Insulator

K. Itoh¹, H. Nakaya¹, Y. Kawakami¹, T. Fukatsu¹, H. Itoh^{1,2}, T. Sasaki^{3,2}, S. Saito⁴, and S. Iwai^{1,2}

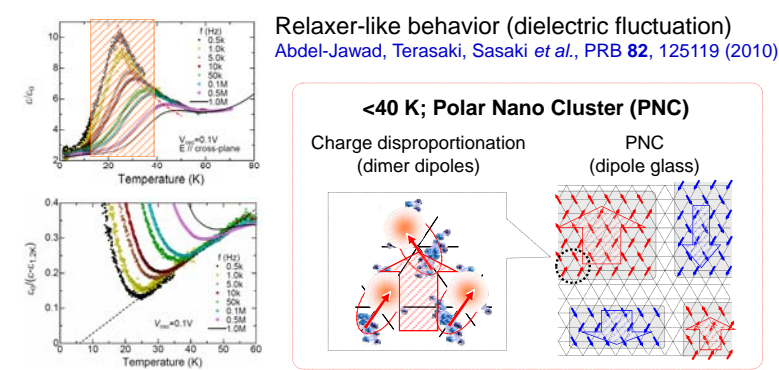
Department of Physics, Tohoku University¹, JST-CREST²,

Institute for Materials Research, Tohoku University³, National Institute of Information and Communications Technology⁴

Low-Dimensional Organic Conductor $\kappa\text{-(ET)}_2\text{Cu}_2(\text{CN})_3$



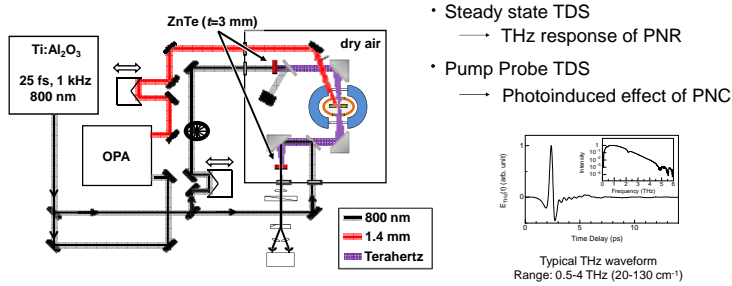
Dielectric anomaly in $\kappa\text{-(ET)}_2\text{Cu}_2(\text{CN})_3$



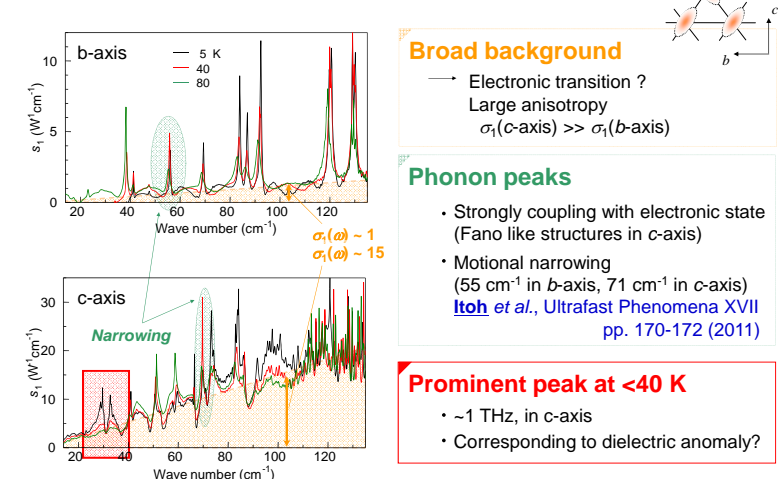
Objective: Capture the optical response of PNC

- Steady state Terahertz (THz) spectrum
- Optical pump THz probe measurements

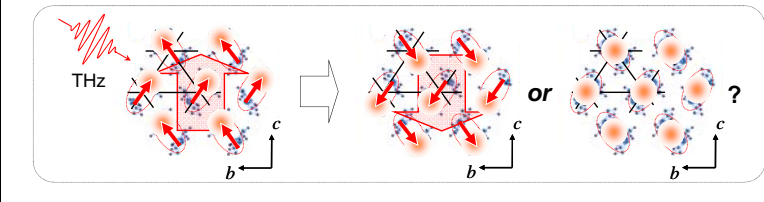
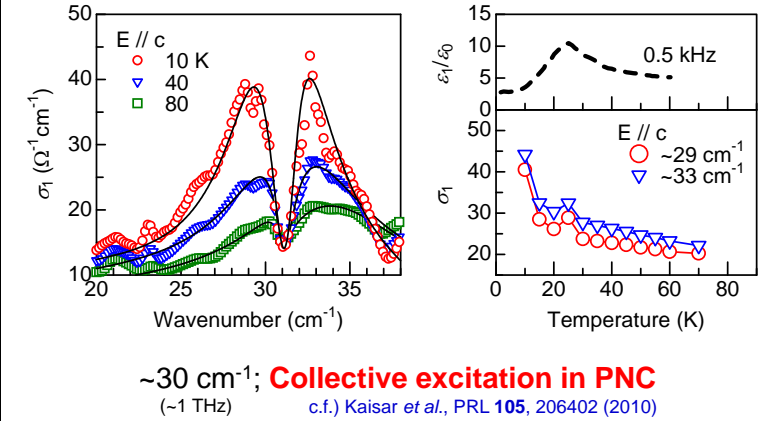
Terahertz Time-Domain Spectroscopy (THz-TDS)



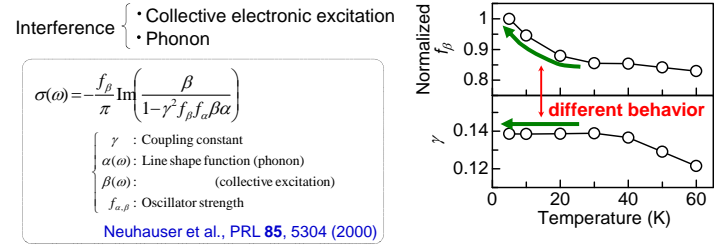
Optical Conductivity in the bc -plane



Optical Transition of Polar Nano Cluster

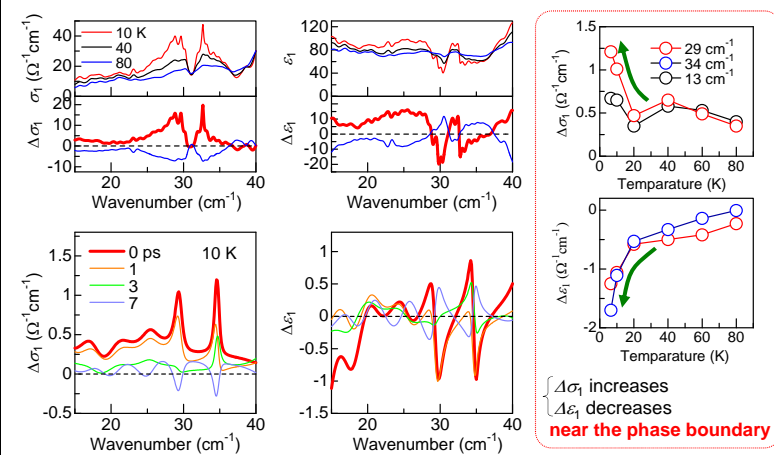


Fano resonance



Photoinduced Proliferation of Polar Nano Cluster

Pump wavelength; 1.4 μm (intra-dimer excitation)



$$\text{Photoinduced } \begin{cases} \Delta\sigma_1 \\ \Delta\epsilon_1 \end{cases} \cong \begin{cases} \sigma_1(10\text{ K}) - \sigma_1(40\text{ K}) \\ \epsilon_1(10\text{ K}) - \epsilon_1(40\text{ K}) \end{cases}$$

Optical Growth of Polar Nano Cluster

Proliferation is driven by the collapse of dimer Mott phase

Summary;

THz response associated with polar nano clusters or dipole glass state was investigated in organic dimer Mott insulator $\kappa\text{-(ET)}_2\text{Cu}_2(\text{CN})_3$ exhibiting relaxer-like low-frequency dielectric anomaly. Electronic transition at ~1 THz can be assigned to the collective excitation of the intradimer dipole in the polar nano cluster. Optical pump and THz probe measurement shows that photoinduced increase of this THz response reflecting the proliferation of the polar cluster. Such photoinduced proliferation of polar cluster is discussed in terms of the competition between the dimer Mott phase and the charge ordered ferroelectric phase.

Let $\mathcal{E}(\mathbb{C})$ and $\mathcal{M}(U)$ denote, respectively, the sets of all entire functions and all meromorphic functions on $U \subseteq \mathbb{C}$. Let Δ_c be the forward difference operator: $\Delta_c f(z) := f(z+c) - f(z)$. The following two results are fundamental in this study.

Theorem 1. (Leont'ev) *Let D be a convex polygonal domain with vertices $\gamma_1, \gamma_2, \dots, \gamma_n \in \mathbb{C}$ and let $\gamma_{n+1} = \gamma_1$. For $1 \leq k \leq n$, let $c_k = \gamma_{k+1} - \gamma_k$ and let D_k be the open half-plane containing D bounded by the line through γ_k and γ_{k+1} . Then, any function f holomorphic in D has a periodic decomposition of the form*

$$f(z) = P_{c_1}(z) + P_{c_2}(z) + \dots + P_{c_n}(z) \quad (z \in D),$$

where each P_{c_k} is holomorphic and c_k -periodic in D_k .

Theorem 2. *Let $c_1, c_2, \dots, c_n \in \mathbb{C}$ be pairwise linearly independent over \mathbb{R} . Then any function $f \in \mathcal{M}(\mathbb{C})$ satisfying the difference equation $\Delta_{c_1} \Delta_{c_2} \dots \Delta_{c_n} f = 0$ has a meromorphic periodic decomposition of the form $f(z) = P_{c_1}(z) + P_{c_2}(z) + \dots + P_{c_n}(z)$, where each $P_{c_k}(z) \in \mathcal{M}(\mathbb{C})$ is c_k -periodic.*

As an application of these theorems, we can show the following theorem.

Theorem 3. *The definitions of D , D_k , γ_k , and c_k are same as in Theorem 1. Let $c_1, c_2, \dots, c_n \in \mathbb{C}$ be pairwise linearly independent over \mathbb{R} (this assumption implies that D has no parallel sides). Then for any function f holomorphic in \overline{D} , the following periodic decomposition holds for any $z \in \overline{D} \setminus \{\text{poles}\}$:*

$$f(z) = P_{c_1}(z) + P_{c_2}(z) + \dots + P_{c_n}(z),$$

where each P_{c_k} is meromorphic and c_k -periodic in D_k . Furthermore if n is even, P_{c_k} and its first $\frac{n}{2} - 1$ derivatives are continuous in $\partial \overline{D}$; and if n is odd, P_{c_k} and its first $\frac{n-1}{2} - 1$ derivatives are continuous in $\partial \overline{D}$.

This result is extendable to the periodic decompositions of meromorphic functions on \overline{D} . Let $f \in \mathcal{M}(\overline{D})$. If there are c_k -periodic meromorphic functions $Q_{c_k} \in \mathcal{M}(\mathbb{C})$ ($1 \leq k \leq n$) such that $f - (Q_{c_1} + Q_{c_2} + \dots + Q_{c_n})$ is holomorphic in \overline{D} , then Theorem 3 gives a meromorphic periodic decomposition $f - (Q_{c_1} + Q_{c_2} + \dots + Q_{c_n}) = P_{c_1} + P_{c_2} + \dots + P_{c_n}$. Thus we have a meromorphic periodic decomposition

$$f = (P_{c_1} + Q_{c_1}) + (P_{c_2} + Q_{c_2}) + \dots + (P_{c_n} + Q_{c_n}).$$

Let $R_{c_k} = P_{c_k} + Q_{c_k}$. Each R_{c_k} is meromorphic in D_k . If we can choose Q_{c_k} as a function that has no poles in $\partial \overline{D}_k$, R_{c_k} is also a c_k -periodic meromorphic function in \overline{D}_k that is holomorphic on $\partial \overline{D}_k$. Then what are conditions on n for all meromorphic functions on \overline{D} to have such decompositions? We give some results for $n = 3, 4, 5, 6$ as follows.

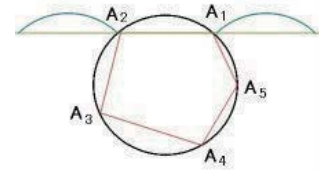
Proposition 1. *We assume $n = 3$, that is to say, D is a triangle domain. Then any meromorphic function f on \overline{D} has a meromorphic periodic decomposition $f(z) = P_{c_1}(z) + P_{c_2}(z) + P_{c_3}(z)$ ($z \in D \setminus \{\text{poles}\}$), where each P_{c_k} is a c_k -periodic meromorphic function on D_k .*

Theorem 4. *We assume $n = 4, 5, 6$. That is to say, D is a convex quadrilateral, pentagonal, or hexagonal domain. Then any meromorphic function f on \overline{D} has a meromorphic periodic decomposition $f(z) = P_{c_1}(z) + P_{c_2}(z) + \dots + P_{c_n}(z)$, where for each P_{c_k} is meromorphic and c_k -periodic in D_k and (i) ($n = 4$) continuous in $\partial \overline{D}_k$, and P_{c_k}' is continuous in $\partial \overline{D}_k$; (ii) ($n = 5$) continuous in $\partial \overline{D}_k$; or (iii) ($n = 6$) continuous in $\partial \overline{D}_k$ except at poles, and its first two derivatives are also continuous in $\partial \overline{D}_k$ except at poles.*

We can consider another application of Theorem 1. We give periodic decompositions of functions holomorphic in domains satisfying some conditions. Let D be a domain and let $\gamma_1 (= \gamma_{n+1}), \gamma_2, \dots, \gamma_n \in \partial D$ and $c_k = \gamma_{k+1} - \gamma_k$. Consider the following conditions:

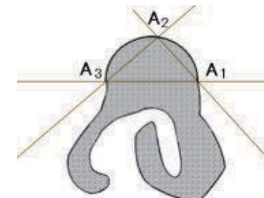
- (a) The n -sided polygon $[\gamma_1 \gamma_2 \dots \gamma_n]$ is convex and contained in D . Let E be this polygonal domain.
- (b) For $1 \leq k \leq n$, let E_k be the open half-plane containing E bounded by the line through γ_k and γ_{k+1} and let $G_k = (\mathbb{C} \setminus \overline{E_k}) \cap E_{k-1} \cap E_{k+1} \cap D$. Then $G_k \cap (G_k + \gamma_k) = G_k \cap (G_k - \gamma_k) = \emptyset$.
- (c) $\partial G_k \cap \partial E_{k-1} = \{\gamma_k\}$, $\partial G_k \cap \partial E_{k+1} = \{\gamma_{k+1}\}$ ($1 \leq k \leq n$).

If the above conditions are satisfied, we call $D = G_1 \cup G_2 \cup \dots \cup \overline{E} \setminus \{\gamma_1, \dots, \gamma_n\}$ a *strictly convex polygonal decomposition* of domain D .



Theorem 5. *Let D be a domain that has a strictly convex polygonal decomposition as above and let $D_k = E_k \cup \{\bigcup_{m \in \mathbb{Z}} (D + m\gamma_k)\}$ ($1 \leq k \leq n$). Then any function f holomorphic in D has a periodic decomposition $f = P_{c_1} + P_{c_2} + \dots + P_{c_n}$, where each P_{c_k} is holomorphic and c_k -periodic in D_k .*

We consider the condition (d) instead of (b), (c):
 (d) $D \cap (\mathbb{C} \setminus \overline{E_k}) \cap \{(\mathbb{C} \setminus \overline{E_{k-1}}) \cup (\mathbb{C} \setminus \overline{E_{k+1}})\} = \emptyset$ ($1 \leq k \leq n$).
 If the conditions (a), (d) are satisfied, we call $D = G_1 \cup G_2 \cup \dots \cup \overline{E} \setminus \{\gamma_1, \dots, \gamma_n\}$ a *convex polygonal decomposition* of domain D . We can think that $G_k + m\gamma_k$ is contained in separate complex planes for each $m \in \mathbb{Z}$. Then D_k is regarded as a Riemann surface that have a c_k -periodic shape and satisfy $G_k \cap (G_k + \gamma_k) = G_k \cap (G_k - \gamma_k) = \emptyset$. Then Theorem 5 can be extended to functions holomorphic in domains having convex polygonal decompositions.



An ODE-diffusion system modeling regeneration of Hydra

Madoka Nakayama

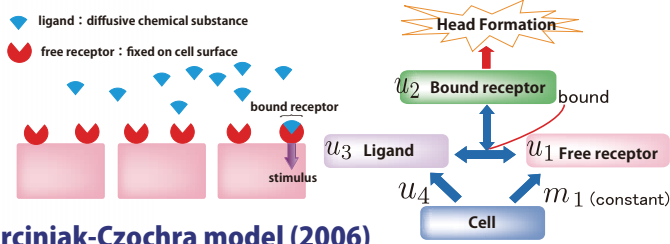
(Mathematical Institute, Tohoku University)

What controls a head regeneration of Hydra?



- **Activator-inhibitor reaction**
Gierer-Meinhardt (1972), MacWilliams (1982)
- **Receptor-Ligand reaction**
Sherrat, Maini, Jäger & Müller (1995),
Marciniak-Czochra (2003, 2006, 2010)

The idea of a Receptor-Ligand reaction



Marciniak-Czochra model (2006)

$$\begin{cases} \tau \frac{\partial u_1}{\partial t} = -\mu_1 u_1 - b u_1 u_3 + d u_2 + m_1 & \text{for } x \in \Omega, t > 0, \\ \tau \frac{\partial u_2}{\partial t} = -\mu_2 u_2 + b u_1 u_3 - d u_2 & \text{for } x \in \Omega, t > 0, \\ \frac{\partial u_3}{\partial t} = \frac{1}{\gamma} \Delta u_3 - \mu_3 u_3 - b u_1 u_3 + d u_2 + u_4 & \text{for } x \in \Omega, t > 0, \\ \frac{\partial u_4}{\partial t} = -\delta u_4 + \frac{m_4 u_3}{1 + \sigma u_4^2 - \beta_l u_4} & \text{for } x \in \Omega, t > 0, \\ u_j(x, 0) = \phi_j(x) \quad (j = 1, \dots, 4) & \text{for } x \in \Omega. \end{cases}$$

u_1 : Density of free receptor,
 u_2 : Density of bound receptor,
 u_3 : Density of ligand,
 u_4 : Production rate of ligand,
 $\tau, \mu_1, \mu_2, \mu_3, m_1, m_4, b, d, \sigma, \alpha, \beta_l, \delta$: positive constants,
 $\partial u_3 / \partial \nu = 0, \phi_j(x) \ (j = 1, \dots, 4)$ are assumed to be smooth and positive on $\bar{\Omega}$.

Boundedness of a solution

There exist positive constants $\rho_j \ (j = 1, \dots, 4)$, depending on the initial functions $(\phi_1, \phi_2, \phi_3, \phi_4)$ s.t.

$$0 < u_j(x, t) \leq \rho_j \quad \text{for all } x \in \Omega, t \geq 0,$$

for $j = 1, 2, 3, 4$.

Diffusion-ODE system

$$\tau = 0$$

$$u(x, t) = u_3(x, t), \quad v(x, t) = u_4(x, t)$$

$$f(u, v) = v - \mu_3 u - \frac{m_1 \mu_2 b u}{\mu_1 (\mu_2 + d) + \mu_2 b u}, \quad g(u, v) = -\delta v + \frac{m_4 u}{1 + \sigma v^2 - \beta_l v}.$$

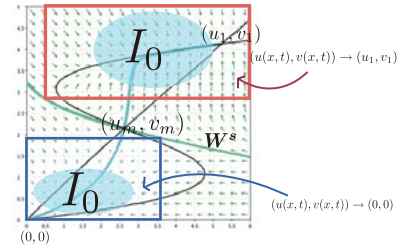
$$\begin{cases} \frac{\partial u}{\partial t} = \frac{1}{\gamma} \frac{\partial^2 u}{\partial x^2} + f(u, v), & \frac{\partial v}{\partial t} = g(u, v) & \text{for } 0 < x < 1, t > 0, \\ \frac{\partial u}{\partial x} = 0 & \text{at } x = 0, 1, t > 0, \\ u(x, 0) = \phi_3(x), v(x, 0) = \phi_4(x) & \text{for } 0 \leq x \leq 1 \end{cases}$$

Constant stationary solutions

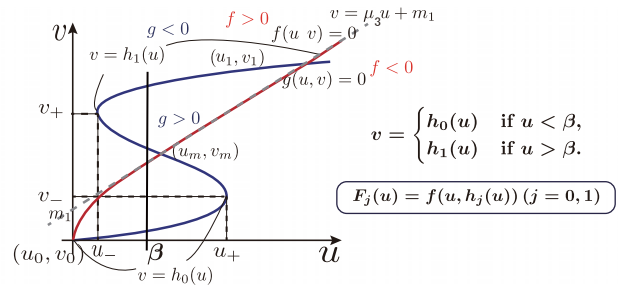
Theorem

Put $I_0 = \{(\phi_3(x), \phi_4(x)) | 0 \leq x \leq 1\}$.

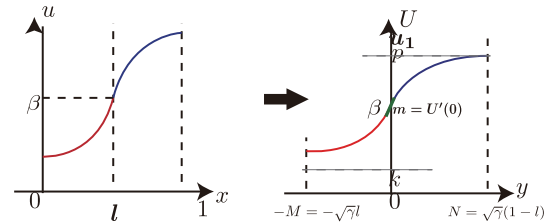
- $\exists (U_0, V_0) \in W^*$ s.t. $0 < U_0 < U_s$ and $I_0 \subset (U_0, \infty) \times (V_0, \infty)$.
Then $(u(x, t), v(x, t)) \rightarrow (u_1, v_1)$ as $t \rightarrow +\infty$.
- $\exists (U_0, V_0) \in W^*$ s.t. $0 < U_0 < U_s$ and $I_0 \subset (0, U_0) \times (0, V_0)$.
Then $(u(x, t), v(x, t)) \rightarrow (u_0, v_0) = (0, 0)$ as $t \rightarrow \infty$.



Monotone increasing stationary solutions



$$(S) \begin{cases} \frac{d^2 u}{dx^2} + \gamma F_0(u) = 0 & \text{for } 0 < x < l, \\ \frac{d^2 u}{dx^2} + \gamma F_1(u) = 0 & \text{for } l < x < 1, \\ u_x(0) = u_x(1) = 0, & u(l) = \beta. \end{cases}$$



$$\mathfrak{F}_0(u) = \int_{u_0}^u F_0(z) dz, \quad \mathfrak{F}_1(u) = \int_{u_1}^u F_1(z) dz$$

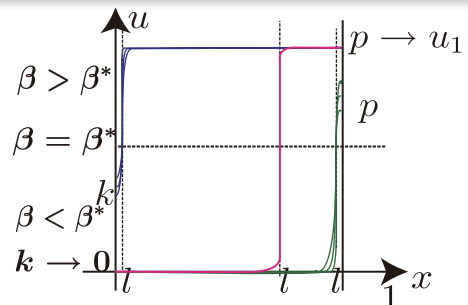
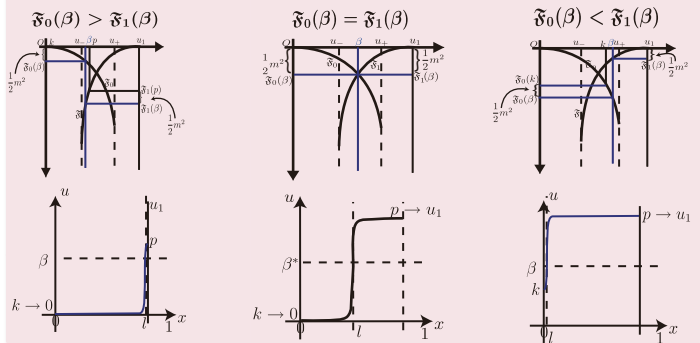
$\forall \beta \in (u_-, \min\{u_+, u_1\})$. \exists one-parameter family of solutions

$$(\gamma(\beta, m), l(\beta, m), u(x; \beta, m)) \in (0, +\infty) \times (0, 1) \times C^1[0, 1]$$

$$\text{with } m \in (0, \min\{\sqrt{2|\mathfrak{F}_0(\beta)|}, \sqrt{2|\mathfrak{F}_1(\beta)|}\})$$

such that

- $u(l(\beta, m); \beta, m) = \beta, \cdot u'(x; \beta, m) > 0 \ (0 < x < 1)$,
- $u(x; \beta, m)$ solves (S) with $\gamma = \gamma(\beta, m)$.



Network Games with and without Synchronicity

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Introduction

In this presentation, we study **mixed Nash equilibria** for stochastic strategies in the **network games**. We then generalize our network game to an **asynchronous game**, where two players repeatedly execute simultaneous games.

Network Games and Profits

Definition 1. Let $G = (V, E)$ be an undirected graph with no isolated vertices. Fix integer α and δ with $\alpha, \delta \geq 1$. A network game $\Gamma_{\alpha, \delta}(G) = \langle \mathcal{N}, \mathcal{S} \rangle$ on G is defined as follows:

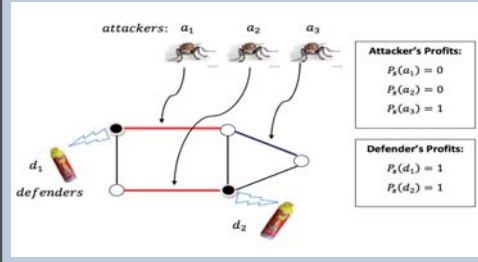
$\mathcal{N} = \mathcal{N}_A \cup \mathcal{N}_D$ is the set of players where \mathcal{N}_A is a finite set of **attackers** a_i , $1 \leq i \leq \alpha$
 \mathcal{N}_D is a finite set of **defenders** d_j , $1 \leq j \leq \delta$
 $\mathcal{S} = E^\alpha \times V^\delta$ is the strategy set of $\Gamma_{\alpha, \delta}(G)$

The individual **profit** of attacker a_i , is given by

$$P_s(a_i) = \begin{cases} 0 & \text{if } v_j \in e_i \text{ for some } j, 1 \leq j \leq \delta \\ 1 & \text{if } v_j \notin e_i \text{ for all } j, 1 \leq j \leq \delta \end{cases}$$

The individual **profit** of defender d_j is given by

$$P_s(d_j) = |\{i : 1 \leq i \leq \alpha, v_j \in e_i\}|$$



Characterizations of Mixed Nash equilibria

A **mixed strategy** for an attacker (resp. defender) is a probability distribution over edges (resp. vertices) of G .

$\sigma_i(e)$: the probability that attacker a_i chooses edge e

$\tau_j(v)$: the probability that defender d_j chooses vertex v .

Definition 3. For each $i \leq \alpha$, $\sigma_i : E \rightarrow [0, 1]$ satisfies $\sum_{e \in E} \sigma_i(e) = 1$, and for each $j \leq \delta$, $\tau_j : V \rightarrow [0, 1]$ satisfies $\sum_{v \in V} \tau_j(v) = 1$.

A **mixed profile** $s = \langle \sigma_1, \dots, \sigma_\alpha, \tau_1, \dots, \tau_\delta \rangle$ is a collection of mixed strategies, one for each player.

Definition 4. The **support** of a player $r \in \mathcal{N}$ in a profile s , denoted by $S_s(r)$, is the set of edges or vertices to which r assigns positive probability in s . Let $S_s(A) = \bigcup_{a_i \in \mathcal{N}_A} S_s(a_i)$ and $S_s(D) = \bigcup_{d_j \in \mathcal{N}_D} S_s(d_j)$.

Let **Save(e)** := the event that at least one end $v \in e$ is protected by a defender.

Let $\pi_s(\text{Save}(e))$ be the probability of the event **Save(e)** with respect to s .

Definition 5 (Expected Profits). For a defender $d_j \in \mathcal{N}_D$,

$$P_s(d_j) = \sum_{\substack{v \in V \\ i \leq \alpha \\ e \in E(v)}} \tau_j(v) \sigma_i(e) = \sum_{v \in V} \tau_j(v) \sum_{\substack{i \\ e \in E(v)}} \sigma_i(e)$$

For an attacker $a_i \in \mathcal{N}_A$,

$$P_s(a_i) = \sum_{e \in E} \sigma_i(e) \cdot (1 - \pi_s(\text{Save}(e)))$$

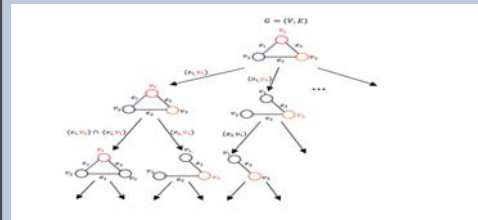
Definition 6 (Nash Equilibrium). A mixed profile s is a **mixed Nash equilibrium** if for each player $r \in \mathcal{N}$, it maximizes P_s over all profiles \bar{s} that differ from s only with respect to the mixed strategy of player r .

Intuitively, no player can gain more by a unilateral change of his strategy. We proceed to study the characterization of a mixed Nash equilibrium.

Theorem 1. In any mixed Nash equilibrium s of $\Gamma_{\alpha, \delta}(G)$, $S_s(D)$ is a vertex cover of G .

Theorem 2. In any mixed Nash equilibrium s of $\Gamma_{\alpha, \delta}(G)$, $S_s(A)$ is an edge cover of the subgraph of G obtained by restricting to $S_s(D)$.

Asynchronous Games



Let $G = (V, E)$ be an undirected graph. The set $\mathcal{P} \subset (E \times V)^{\mathbb{N}}$ of partial plays is defined recursively as follows:
 -Put the empty sequence $\lambda \in \mathcal{P}$ and $E_\lambda = E$ and $V_\lambda = V$.
 -Now assume that $\eta \in \mathcal{P}$, and $E_\eta \subset E$ and $V_\eta \subset V$ have been defined. We put $\rho := \eta \frown \langle (e, v) \rangle$ into \mathcal{P} , if $e \in E_\eta$ and $v \in V_\eta$. Then, we define

$$E_\rho := \begin{cases} E_\eta & \text{if } v \in e \\ E_\eta - \{e\} & \text{if } v \notin e \end{cases}$$

$$V_\rho := V(E_\rho).$$

Finally, let $[\mathcal{P}] = \{w \in (E \times V)^{\mathbb{N}} : \text{each finite initial segment of } w \text{ belongs to } \mathcal{P}\}$.

Determinacy of Asynchronous Games

Theorem 3. Given a graph G and function $f : \{G' \subsetneq G : G' \text{ a subgraph}\} \rightarrow \mathbb{R}$. Then, the one-round game $\Gamma(G; f)$ is determined with a stable solution.

Proof. Suppose $f : \{G' \subsetneq G\} \rightarrow \mathbb{R}$, and $f(G) = x \in \mathbb{R}$ (for the next round with the same G) are given. The expected profit according to a profile (σ, τ) is following:

$$P(\sigma, \tau, x) := \sum_{v \notin e} \sigma(e) \tau(v) \{1 + f(h(e, v))\} + \sum_{v \in e} \sigma(e) \tau(v) \frac{1}{2} x.$$

By the mini-max theorem, we have $\max_\sigma \min_\tau P(\sigma, \tau, x) = \min_\tau \max_\sigma P(\sigma, \tau, x)$ for all $x \in \mathbb{R}$. Note that the set of strategy σ 's (similar for τ 's) is bounded closed convex, and so it is easy to see that $M(x) = \max_\sigma \min_\tau P(\sigma, \tau, x)$ is a continuous function.

Now put $m = 1 + \max f$. Clearly, $M : [0, m] \rightarrow [0, m]$. So it has a fixed point \hat{x} . Then $M(\hat{x})$ serves as a stable solution. \square

Now, we define a function $f : W \rightarrow \mathbb{R}$, which plays the role of a natural payoff for the attacker of our asynchronous game. For $w = ((e^1, v^1), (e^2, v^2), (e^3, v^3), \dots) \in [\mathcal{P}]$, we set

$$b_i(w) := \begin{cases} 0 & \text{if } v^i \in e^i, \\ 1 & \text{if } v^i \notin e^i. \end{cases}$$

Then we define

$$f(w) = \sum_{i > 0} \frac{b_i(w)}{2^i}$$

Theorem 4. The asynchronous game $\Gamma(G; f)$ is determined.

Proof. (Idea) We show the existence of a stable solution in the game without referring to the Blackwell determinacy and we reduced the infinite game $\Gamma(G; f)$ to a finite-round game. \square

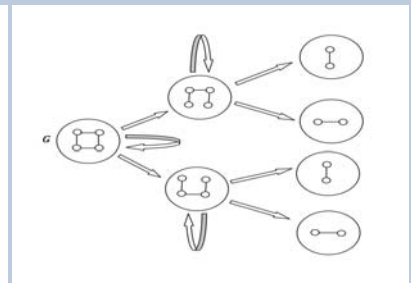
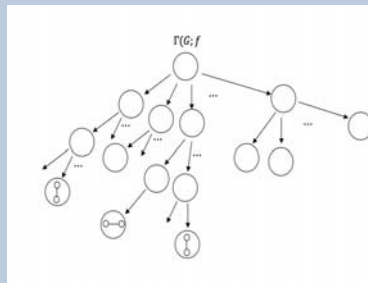
Expected Profits (One-Round)

Given $G :=$ a finite graph and $f : \{G' : G' \text{ is a proper subgraph of } G\} \rightarrow \mathbb{R}$, a state evaluation. We define a function $h : E \times V \rightarrow \{\text{subgraphs of } G\}$ by

$$h(e, v) := \begin{cases} G & \text{if } v \in e \\ G(E \setminus \{e\}) & \text{if } v \notin e \end{cases}$$

Definition 2. Let $s = (\sigma, \tau)$ be a pair of mixed strategies. The **expected profit** of the attacker with the delay constant c given by

$$P_s(\Gamma(G; f)) := \sum_{v \notin e} \sigma(e) \tau(v) \{1 + f(h(e, v))\} + \sum_{v \in e} \sigma(e) \tau(v) c P_s(\Gamma(G; f))$$



Asymptotic Behavior of Non-local Feynman-Kac Semigroups

Masakuni MATSUURA*

February 21, 2012

1 Feynman-Kac penalization problem

Let X be a symmetric α -stable process on \mathbb{R}^n ($n > \alpha$) and let A_t^μ be a positive continuous additive functional corresponding with Revuz measure μ . We define $A_t^{\mu,F}$ by

$$A_t^{\mu,F} := A_t^\mu + \sum_{u \leq t} F(X_{u-}, X_u),$$

where $F(x, y)$ is a bounded positive symmetric function vanishing on the diagonal. We call the next problems the Feynman-Kac penalization problem.

(1) Does there exist a probability measure $\tilde{\mathbb{P}}_x$ such that

$$\lim_{t \rightarrow \infty} \frac{\mathbb{E}_x[e^{A_t^{\mu,F}} \mathbf{1}_\Lambda]}{\mathbb{E}_x[e^{A_t^{\mu,F}}]} = \tilde{\mathbb{P}}_x[\Lambda] \text{ for every } \Lambda \in \mathcal{F}_s?$$

(2) Is the limit distribution $\tilde{\mathbb{P}}_x$ determined by a martingale M_s : $d\tilde{\mathbb{P}}_x = M_s d\mathbb{P}_x$?

We first decompose the non-local multiplicative functional $e^{A_t^{\mu,F}}$ as $L_t e^{A_t^{\mu+\mu_{F_1}}}$, where

$$L_t := e^{\sum_{u \leq t} F(X_{u-}, X_u) - c \int_0^t \int F_1(X_u, y) |X_u - y|^{-(n+\alpha)} dy du}$$

is a martingale and

$$\mu_{F_1}(dx) := c \left(\int F_1(x, y) |x - y|^{-(n+\alpha)} dy \right) dx, F_1 = e^F - 1.$$

$c > 0$ is a constant. We also transform X to Y whose law is $dP_x^L = L_t dP_x$. We then find that the Dirichlet form \mathcal{E}^Y of Y is

$$\mathcal{E}^Y(u, u) = \frac{c}{2} \int (u(x) - u(y))^2 e^{F(x, y)} |x - y|^{-(\alpha+n)} dx dy.$$

Let $\lambda_0 := \inf\{\mathcal{E}^Y(u, u); \|u\|_{L^2(\mu+\mu_{F_1})} = 1\}$. Feynman-Kac penalization problem is solved for $\lambda_0 \neq 1$ if μ and μ_{F_1} are in the Green-tight Kato class \mathcal{K}_∞ . (a) If $\lambda_0 > 1$, then M_s is given by

$$M_s = \frac{e^{A_s^{\mu,F}} h(X_s)}{h(x)}, \quad h(x) := \mathbb{E}_x^M[e^{A_\infty^{\mu+\mu_{F_1}}}]$$

(b) If $\lambda_0 < 1$, then M_s is given by

$$M_s = \frac{e^{-\theta_0 s + A_s^{\mu,F}} h(X_s)}{h(x)},$$

where $\theta_0 > 0$ is the constant such that $\inf\{\mathcal{E}_{\theta_0}^Y(u, u); \|u\|_{L^2(\mu+\mu_{F_1})} = 1\} = 1$. h is the ground

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state of Schrödinger type operator $\mathcal{L}^Y - (\mu + \mu_{F_1})$. Here, \mathcal{L}^Y is the generator of Y : $\mathcal{E}^Y(u, u) = (\mathcal{L}^Y u, u)$.

(c) If $\lambda_0 = 1$, then M_s is given by

$$M_s = \frac{e^{A_s^{\mu,F}} h(X_s)}{h(x)}$$

only if μ and μ_{F_1} are in the special Kato class $\mathcal{K}_s(\subset \mathcal{K}_\infty)$. Here, h is the ground state as in the case of (b).

2 Examples of jumping functions

We give examples of jumping functions with full support such that $\mu_F \in \mathcal{K}_\infty$

$$F(x, y) := (1 \wedge |x - y|^p) \langle x \rangle^{-q} \langle y \rangle^{-q} \text{ for } p > \alpha \text{ and } q > n.$$

Here, $\langle \cdot \rangle := \sqrt{1 + |\cdot|^2}$. If we further assume $q > 2n - \alpha$, then $\mu_F \in \mathcal{K}_s$.

3 Asymptotic behavior of non-local Feynman-Kac semigroups

We also see asymptotic behavior of non-local Feynman-Kac semigroups. Let

$$P_t^{\mu,F} f(x) := \mathbb{E}_x[e^{A_t^{\mu,F}} f(X_t)], \quad \text{for } f \in L^2 \cap \mathcal{B}_b.$$

We then obtain

$$P_t^{\mu,F} \mathbf{1}(x) \begin{cases} \sim \mathbb{E}_x[e^{A_\infty^{\mu,F}}] & \text{if } \lambda_0 > 1 \\ \sim (h(x) \int_{\mathbb{R}^n} h(x) dx) e^{\theta_0 t} & \text{if } \lambda_0 < 1 \\ \sim (h(x) \int_{\mathbb{R}^n} h(x) d(\mu + \mu_{F_1})) t & \text{if } \lambda_0 = 1, n > 2\alpha \\ = o(t) & \text{if } \lambda_0 = 1, \alpha < n \leq 2\alpha \end{cases}$$

as $t \rightarrow \infty$.

We further compute the growth of L^p -spectral bounds due to Tawara ([2, Example 5.3]). Let $l_p := \lim_{t \rightarrow \infty} (-1/t) \log \|P_t^{\mu,F}\|_{p,p}$ ($1 \leq p \leq \infty$). We then see

$$l_p = \begin{cases} 0 & \text{if } \lambda_0 \geq 1 \\ -\theta_0 & \text{if } \lambda_0 < 1. \end{cases}$$

References

- [1] Matsuura, M., *Feynman-Kac Penalization Problem for Additive Functionals with Jumping Functions*, in preparation.
- [2] Tawara, Y., *L^p -independence of Growth Bounds of Generalized Feynman-Kac Semigroups*, Doctor Thesis, Tohoku University, 2008.

The characterization of a pinned polymer

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Introduction

We consider a polymer which given in the form of the Brownian motion (B_t, P_0) . Let $L_{0,t} = \int_0^t \delta_0(B_s) ds$ denote the energy of the polymer and set

$$\begin{aligned} P_{\beta,t}(d\omega) &= Z_{0,t,\beta}^{-1} \exp(\beta L_{0,t}) P_0(d\omega), \\ Z_{0,t,\beta} &= E_0[\exp(\beta L_{0,t})]. \end{aligned} \quad (1)$$

Where a parameter $\beta > 0$ means the inverse temperature. K. S. Alexander and V. Sidoravicius[1] introduced the *pinned polymer*. That is, for given β , there exists $\delta > 0$ such that $\lim_{t \rightarrow \infty} P_{\beta,t}(L_t/t > \delta) = 1$. Our aim is to give the condition the pinning occurs under the general situation.

1 Preliminaries

Let $M^\alpha = (P_x, X_t)$ be a symmetric α -stable process on \mathbb{R}^d with $0 < \alpha \leq 2$ and $d \leq 2\alpha$ (in particular, the case of $\alpha = 2$ means the Brownian motions). For a positive Radon measure μ belongs to Green tight Kato class (in notation $\mu \in \mathcal{K}_{d,\alpha}^\infty$), let A_t^μ be the positive continuous additive functional under the Revus correspondence. Where the above functional A_t^μ is a extension of $L_{0,t}$. And we define a Gibbs measure $P_{\beta,t}^\mu$ same as (1) by the A_t^μ and the law of α -stable process.

From the large deviation principle, we can estimate the asymptotic behavior of A_t^ν/t : for any $\kappa > 0$,

$$P_{\beta,t}^\mu \left(\frac{A_t^\nu}{t} > \kappa \right) \sim \exp \left\{ -t \inf_{z > \kappa} C_{\beta,N}^*(z) \right\}, \quad \text{as } t \rightarrow \infty. \quad (2)$$

Where a function $C_{\beta,N}^*$ is a Fenchel-Legendre transform of $C_{\beta,N}(\lambda)$: $C_{\beta,N}^*(z) = \sup\{\lambda z - C_{\beta,N}(\lambda) | \lambda \in \mathbb{R}\}$

$$\begin{aligned} C_\beta(\lambda) &:= \lim_{t \rightarrow \infty} \frac{1}{t} \log E_x[\exp\{\beta A_t^\mu + \lambda A_t^\nu\}], \\ C_{\beta,N}(\lambda) &:= C_\beta(\lambda) - C_\beta(0). \end{aligned} \quad (3)$$

We set

$$\begin{aligned} \lambda_0(\beta) &= \sup\{\lambda \in \mathbb{R} | C_\beta(\lambda) = 0\}, \\ \beta_{cr} &= \sup\{\beta > 0 | C_\beta(0) = 0\}, \end{aligned}$$

and

$$z_0(\beta) = \frac{d}{d\lambda} C_\beta(\lambda) \Big|_{\lambda=0}.$$

The function $C_\beta(\lambda)$ is differentiable for every λ . So we can define this. The logarithmic moment generating function $C_\beta(\lambda)$ is called a rate function of the large deviation principle, and it coinside with the spectral bottom of Schrödinger operator $\mathcal{H}_{\mu,\nu}[2]$:

$$\begin{aligned} \mathcal{H}_{\mu,\nu}(\beta, \lambda) &= \frac{1}{2}(-\Delta)^{\frac{\alpha}{2}} - \beta\mu - \lambda\nu, \\ C_\beta(\lambda) &= -\inf\{\sigma(\mathcal{H}_{\mu,\nu}(\beta, \lambda))\}. \end{aligned}$$

In the case of the 1-dimensional Brownian motion and $\mu = \nu = \delta_0$, we can compute the rate function:

$$C_{\beta,N}(\lambda) = \begin{cases} \frac{1}{2}(\lambda^2 + 2\beta\lambda), & \lambda \geq -\beta \\ -\frac{1}{2}\beta^2, & \lambda < -\beta \end{cases}$$

and

$$C_{\beta,N}^*(z) = \begin{cases} \frac{1}{2}(z - \beta)^2, & z > 0 \\ \frac{1}{2}\beta^2, & z = 0 \\ \infty, & z < 0. \end{cases} \quad (4)$$

Therefore, polymer is pinned, for any $\beta > 0$. And $z_0(\beta) = \beta$, $\lambda_0(\beta) = -\beta$, $\beta_{cr} = 0$.

2 Main theorem

In this section, we fix measures $\mu, \nu \in \mathcal{K}_{d,\alpha}^\infty$.

Theorem 2.1. *If $\beta > \beta_{cr}$, then for any $0 < \kappa < z_0(\beta)$,*

$$\lim_{t \rightarrow \infty} P_{\beta,t}^\mu \left(\frac{A_t^\nu}{t} > \kappa \right) = 1. \quad (5)$$

But for $0 < \beta \leq \beta_{cr}$ and any $\kappa > 0$, the limit of (5) is not one.

Skech of the proof.

By the (2), we see that the polymer is pinned if and only if a null set of $C_{\beta,N}^*$ is in \mathbb{R}_+ . Note the null set as $\mathcal{N} = \{z > 0 | C_{\beta,N}^*(z) = 0\}$. In fact, $C_{\beta,N}(\lambda)$ is strictly increasing on $\{\lambda \geq \lambda_0(\beta)\}$ but $C_{\beta,N} \equiv 0$ otherwise, for each $\beta > 0$. This is a key point, and it implies following lemma.

Lemma 2.1. *For each β ,*

$$\begin{aligned} \lambda_0(\beta) < 0 &\iff \mathcal{N} = \{z_0(\beta)\}, \\ \lambda_0(\beta) \geq 0 &\iff \mathcal{N} = \emptyset. \end{aligned}$$

Intuitively, the function $z_0(\beta)$ is increasing by the definition (3). Thus $z_0(\beta) > 0$ if β is large enough. Precisely,

Lemma 2.2.

$$\begin{aligned} 0 < \beta \leq \beta_{cr} &\iff \lambda_0(\beta) \geq 0, \\ \beta > \beta_{cr} &\iff \lambda_0(\beta) < 0. \end{aligned}$$

So we have the conclusion above.

Remark 2.1. *In general, if an underlying process is transient, then $\lambda_0(0) > 0$, otherwise $\lambda_0(0) = 0$. Thus in the recurrent case, $\beta_{cr} = 0$. That is, for any $\beta > 0$, the polymer is pinned when the recurrent case(The d -dimensional Brownian motion is recurrent if $d = 1, 2$, but transient $d \geq 3$). Because α -stable process on \mathbb{R}^2 with $\alpha < 2$ is transient, so we can chose $\beta_{cr} > 0$.*

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- [2] M. Takeda, Asymptotic properties of generalized Feynman-Kac functionals, *Potential Analysis* **9** 261-291, 1998.

On some orbit spaces of prehomogenous vector spaces

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1 What is Prehomogenous vector space?

Let k be an arbitrary field. Let G be a connected reductive group, V a representation of G , and χ a non-trivial rational character of G , all defined over k .

Definition 1.1 (M.Sato). A triple (G, V, χ) is called a **prehomogenous vector space (PV)** if the following two conditions are satisfied:

- (1) There exists a Zariski open orbit.
- (2) There exists a polynomial $\Delta \in k[V]$ such that $\Delta(gx) = \chi(g)^a \Delta(x)$ for some positive integer a .

If V is irreducible, we may use the notation (G, V) instead of (G, V, χ) , because of the choice of χ is essentially unique. Furthermore, the set $V^{\text{ss}} = \{x \in V; \Delta(x) \neq 0\}$ is well-defined (Where “ss” means “semi stable”). **Roughly speaking, a PV is a representation space of an algebraic group with a polynomial, so called relative invariant polynomial.**

Example 1.2. Put $G = \text{GL}(1) \times \text{GL}(2)$ and $V =$ (the space of binary quadratic forms). For $g = (t, g_1) \in G, x = x(v) \in V$, we define $gx = tx(vg_1)$. Then the couple (G, V) is a PV. Now we identify V with k^3 via $k^3 \ni x = (x_0, x_1, x_2) \mapsto F_x(v) = x_0v_1^2 + x_1v_1v_2 + x_2v_2^2 \in V$. Under this identification, gx is given by $F_{gx}(v) = tF_x(vg_1)$ for $g = (t, g_1) \in G, x \in V$. If we denote this representation by $\rho : G \rightarrow \text{GL}(V)$, we have

- (1) $T = \ker \rho = \{(t^{-2}, \begin{pmatrix} t & 0 \\ 0 & t \end{pmatrix}); t \in \text{GL}(1)\}$,
- (2) $P(gx) = \chi(g)P(x)$, where $P(x) = x_1^2 - 4x_0x_2$ and $\chi(g) = t^2(\det g_1)^2$.

2 Parametrization of fields

If (G, V) be a (regular irreducible) PV, then $V_{\bar{k}}$ is a single $G_{\bar{k}}$ -orbit in general. But V_k is not necessarily a single G_k -orbit, and it is difficult and interesting problem to describe the orbit space $G_k \backslash V_k$ explicitly. In this section, we consider an arithmetical meaning of the orbit space of PV's.

In the rest of this section, we will assume (G, V) as in Example 1.2. For $x \in V$, $k(x)$ denotes the splitting field of $F_x(v)$. Note that F_x is separable over k if and only if $x \in V_k^{\text{ss}}$.

Proposition 2.1 (Gauss). For $x, y \in V_k^{\text{ss}}, G_k x = G_k y$ if and only if $k(x) = k(y)$. Therefore, the map

$$\begin{array}{ccc} G_k \backslash V_k^{\text{ss}} & \rightarrow & \{\text{separable extension of } k \text{ with } \deg \leq 2\} \\ \cup & & \cup \\ G_k x & \mapsto & k(x) \end{array}$$

is well-defined and bijective.

Proof. Since “if” part is obvious, we only prove “only if” part.

(i) $k(x) = k$: Put $w = (0, 1, 0)$. Since F_x splits in k and $x \in V_k^{\text{ss}}$, there exist $t, a, b, c, d \in k$ such that

$$F_x(v) = t(av_1 + cv_2)(bv_1 + dv_2), \quad t \neq 0, ad - bc \neq 0.$$

Therefore $x = (t, \begin{pmatrix} a & b \\ c & d \end{pmatrix})w$.

(ii) $[k(x) : k] = 2$: We can assume $x_0 = y_0 = 1$. If F_x and F_y split in k^{sep} as

$$F_x(v) = (v_1 + \alpha_1 v_2)(v_1 + \alpha_2 v_2), \quad F_y(v) = (v_1 + \beta_1 v_2)(v_1 + \beta_2 v_2),$$

then $k(x) = k(\alpha_1), k(y) = k(\beta_1)$. So we have $k(\alpha_1) = k(\beta_1)$. Therefore, there exist $p \in k^\times, q \in k$ such that $\beta_1 = p\alpha_1 + q$. Since $v_1 + \beta_1 v_2 = (v_1 v_2) \begin{pmatrix} 1 \\ \beta_1 \end{pmatrix} = (v_1 v_2) \begin{pmatrix} 1 & 0 \\ q & p \end{pmatrix} \begin{pmatrix} 1 \\ \alpha_1 \end{pmatrix}$ and this conjugation is equal to $(v_1 v_2) \begin{pmatrix} 1 \\ \beta_2 \end{pmatrix} = (v_1 v_2) \begin{pmatrix} 1 & 0 \\ q & p \end{pmatrix} \begin{pmatrix} 1 \\ \alpha_2 \end{pmatrix}$, we have $y = (1, \begin{pmatrix} 1 & 0 \\ q & p \end{pmatrix})x$. This prove the proposition. \square

Next we describe the structure of the stabilizer. For $x \in V_k^{\text{ss}}, G_x^\circ$ denotes the identity component of G_x .

Proposition 2.2. (1) $[G_x : G_x^\circ] = 2$

$$(2) \quad G_x^\circ \cong \begin{cases} \text{GL}(1) \times \text{GL}(1) & \text{if } k(x) = k \\ \mathfrak{R}_{k(x)/k}(\text{GL}(1)) & \text{if } [k(x) : k] = 2 \end{cases}$$

Note that the “tamagawa number” of G_x°/T is constant multiple of $h_{k(x)} R_{k(x)}$ if $[k(x) : k] = 2$.

3 Application

The following type theorem is so called **density theorems**:

Theorem 3.1 (Goldfeld–Hoffstein 1986).

$$\begin{aligned} \lim_{X \rightarrow \infty} \frac{1}{X^{3/2}} \sum_{\substack{[F:\mathbb{Q}]=2 \\ 0 < -D_F < X}} h_F R_F &= \frac{\pi}{36} \prod_p (1 - p^{-2} - p^{-3} + p^{-4}) \\ \lim_{X \rightarrow \infty} \frac{1}{X^{3/2}} \sum_{\substack{[F:\mathbb{Q}]=2 \\ 0 < D_F < X}} h_F R_F &= \frac{\pi^2}{36} \prod_p (1 - p^{-2} - p^{-3} + p^{-4}) \end{aligned}$$

The above theorem was first proved by using Eisenstein series of half-integral weight. In 1993, Datskovsky gave another proof based on the “zeta function” for the space of binary quadratic forms.

References

- [1] A.Yukie, Rational orbit decomposition of prehomogenous vector spaces (Lecture note), Available from <http://www.math.tohoku.ac.jp/~yukie>

Weak Determinacy of Infinite Games and Corresponding Hierarchy of Inductive Definitions, Keisuke Yoshii, (D1)

The PURPOSE of this research is ...

to investigate the logical strength of weak determinacy of infinite game from the standpoint of reverse mathematics.

In Other Words ...

How strong axiom is exactly needed to compute the winning strategy for certain game?

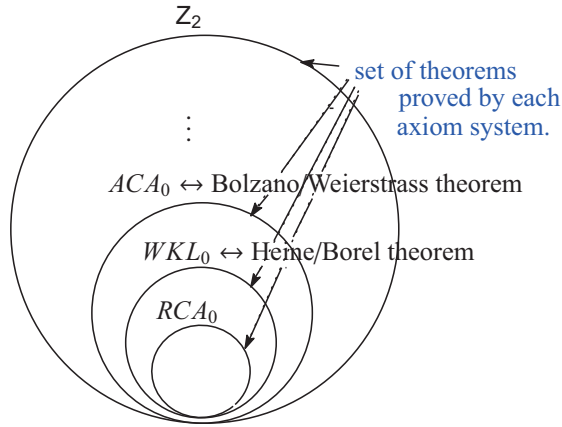
What is the axiom X ?!

$X \vdash$ A winning strategy exists for game G .

• \vdash : “ $T \vdash \varphi$ ” means “ φ is provable from T .”

Reverse Mathematics Program: a program to answer the next questions: *What set existence axioms are needed to prove the theorems of ordinary mathematics?*

Examples of Axioms of Z_2 and equivalent Theorems



- $RCA_0, WKL, ACA, \dots, Z_2$: axiom systems of second order arithmetic, from weaker to stronger.
- Z_2 : strongest axiom system.
- \leftrightarrow : logical equivalence relation.

Determinacy of Infinite Games

- Let $C \subseteq \mathbb{N}^{\mathbb{N}}$. Fix $A \in C$, set of infinite sequences of \mathbb{N} .
- Player I and II alternately choose natural numbers as follows:

I n_0 n_2 ...
 II n_1 ...

- player I **wins** if $n_0, n_1, n_2, \dots \in A$. (II wins o.w.)
- C -Game G_A is **determinate** (det.) if one of the players has a winning strategy for G_A .

Mathematical Meanings of game determinacy?

“A game is determinate” = “Some real number with some complexity exists”!

Indeed, each axiom in Z_2 , g.e. RCA_0, WKL, \dots , also insists

an existence of set of real number with some complexity.

Some well-known facts on Det. of Games

- $Z_2 \not\vdash \text{Borel}(\Delta_1^1)$ game is determinate.
- $ZFC \vdash$ Borel game is determinate.
- $ZF + AD \vdash$ any set of real numbers are Lebesgue measurable.
- $ZFC + AD \vdash 1 = 0$ (contradiction).

• \neg : “ $T \not\vdash \varphi$ ” means “ φ is Not provable from T .”

• **ZFC**: Zermelo/Fraenkle set theory+Axiom of Choice.

• When A is a Borel set, G_A is called **Borel game**.

• **AD** (axiom of determinacy): For any $A \subseteq \mathbb{N}^{\mathbb{N}}$, G_A is determinate.

Weak Determinacy of Games in Z_2

We consider “How strong axioms are needed to prove the det. of the following classes?” e.g. when $A \in \Sigma_1^0$,

$?? \vdash$ Game G_A is determinate, (written as Σ_1^0 -Det)

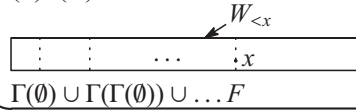
Class on $\mathbb{N}^{\mathbb{N}}$: $\Delta_1^0, \Sigma_1^0, \Sigma_2^0, \Sigma_2^0 \wedge \Pi_2^0, \dots$

(more complex) \rightarrow

Inductive Definitions

Definition 0.1. Σ_1^1 -ID asserts that for any Σ_1^1 -operator Γ , there exists a pre-well-ordering set $W \subset \mathbb{N} \times \mathbb{N}$ on its field F s.t.

- $\forall x \in F \quad W_x = \Gamma(W_{<x}) \cup W_{<x}$,
- $\Gamma(F) \subset F$.



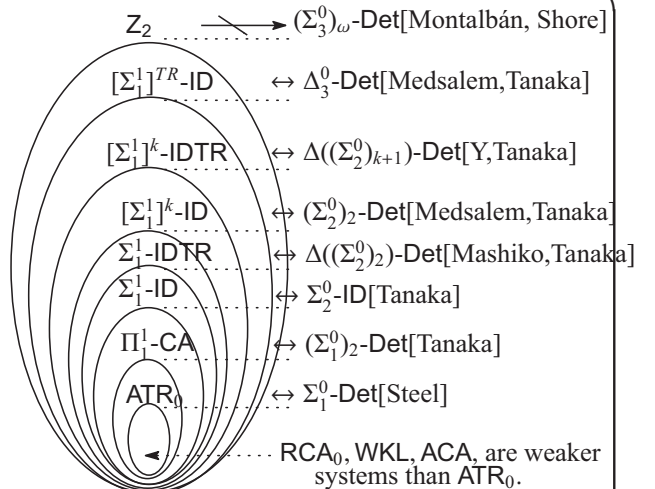
Following variations of inductive definitions have been considered.

Definition 0.2. Σ_1^1 -IDTR:=transfinite iterations of Σ_1^1 -ID.

$[\Sigma]^k$ -ID:= Σ_1^1 -ID with k -operators.

$[\Sigma_1^1]^k$ -IDTR:=trans. iterat. of Σ_1^1 -ID with k -operators.

I.D. v.s. Det. of Games



Location of the concentration point in the ground-state solution of a reaction-diffusion equation in a heterogeneous medium

Hiroko Yamamoto

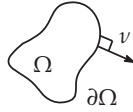
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1 Introduction

Consider the Neumann problem for a semilinear elliptic equation:

$$\begin{cases} \varepsilon^2 \Delta u - a(x)u + b(x)u^p = 0 & \text{in } \Omega, \\ u(x) > 0 & \text{in } \Omega, \\ \frac{\partial u}{\partial \nu} \Big|_{\partial \Omega} = 0, \end{cases} \quad (1)$$

where $\Omega \subset \subset \mathbb{R}^n$, $\partial \Omega \in C^\infty$, and ν denotes the unit outward normal to $\partial \Omega$; $\varepsilon > 0$, $\Delta = \sum_{i=1}^n \partial^2 / \partial x_i^2$; $a, b \in C^\alpha(\bar{\Omega})$; $\min_{x \in \bar{\Omega}} a(x), \min_{x \in \bar{\Omega}} b(x) > 0$. Let the exponent p satisfy $1 < p < (n+2)/(n-2)$ if $n \geq 3$, $1 < p < \infty$ if $n = 1, 2$.

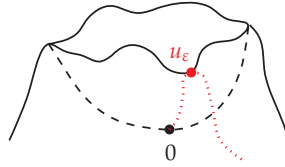


Definition 1. An energy functional of (1) :

$$J_\varepsilon(u) := \frac{1}{2} \int_\Omega \{ \varepsilon^2 |\nabla u(x)|^2 + a(x)u(x)^2 \} dx - \frac{1}{p+1} \int_\Omega b(x)u_+(x)^{p+1} dx$$

where $u_+(x) := \max\{u(x), 0\}$.

By the Mountain Pass Lemma, there exists a critical value c_ε of I_ε . Let $u_\varepsilon \in H^1(\Omega)$ be a critical point for c_ε (i.e. $c_\varepsilon = I_\varepsilon(u_\varepsilon)$). We call u_ε a **ground-state solution** of (1).



Definition 2. A family of ground-state solutions $\{u_\varepsilon\}_{\varepsilon>0}$ **concentrates** at $P_0 \in \bar{\Omega}$.

$\stackrel{\text{def}}{\Leftrightarrow} \exists \{\varepsilon_j\}_{j \in \mathbb{N}} (\varepsilon_j > 0), \exists \{P_j\}_{j \in \mathbb{N}} \subset \bar{\Omega}$ s.t.

- (1) $\varepsilon_j \rightarrow 0, P_j \rightarrow P_0$,
- (2) $\max_{x \in \bar{\Omega}} u_{\varepsilon_j}(x) = u_{\varepsilon_j}(P_j)$,
- (3) $I_{\varepsilon_j}(u_{\varepsilon_j}) = O(\varepsilon_j^n)$.

Known Results	$a(x), b(x)$	Concentration Point
W.-M. Ni & I. Takagi ([1])	$a(x) \equiv b(x) \equiv 1$,	Maximum Point of Mean Curvature of $\partial \Omega$.
X. Ren ([2])	$a \equiv 1$, $b \in C^\alpha(\bar{\Omega}), b(x) > 0$,	<ul style="list-style-type: none"> • $\max_\Omega b(x) > 2^{\frac{p-1}{2}} \max_{\partial \Omega} b(x)$ \Rightarrow Max. Pt. of $b(x)$ in Ω. • $\max_\Omega b(x) < 2^{\frac{p-1}{2}} \max_{\partial \Omega} b(x)$ \Rightarrow Max. Pt. of $b(x)$ on $\partial \Omega$.

2 Main Results

Theorem 3. A function Φ is given by

$$\Phi(x) := a(x)^{\frac{2}{p-1}+1-\frac{n}{2}} b(x)^{-\frac{2}{p-1}}.$$

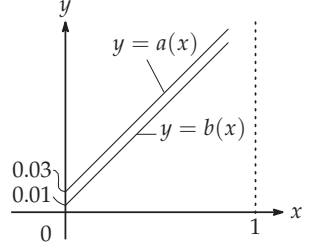
Assume $u_{\varepsilon_j}(P_j) = \max_{x \in \bar{\Omega}} u_{\varepsilon_j}(x)$, $P_j \rightarrow P_0$ as $\varepsilon_j \downarrow 0$. Then:

- (i) $\min_{x \in \bar{\Omega}} \Phi(x) > \frac{1}{2} \min_{x \in \partial \Omega} \Phi(x) \Rightarrow P_0 \in \partial \Omega, \Phi(P_0) = \min_{x \in \bar{\Omega}} \Phi(x)$.
- (ii) $\min_{x \in \bar{\Omega}} \Phi(x) < \frac{1}{2} \min_{x \in \partial \Omega} \Phi(x) \Rightarrow P_0 \in \Omega, \Phi(P_0) = \min_{x \in \bar{\Omega}} \Phi(x)$.

Quiz 4. Where on $[0, 1]$ do you think u_ε concentrates in the following case? ($n = 1, p = 2$)

$$\begin{cases} a(x) = x + 0.03, \\ b(x) = x + 0.01. \end{cases}$$

$$\Rightarrow \Phi(x)^2 = a(x)^5 b(x)^{-4}.$$



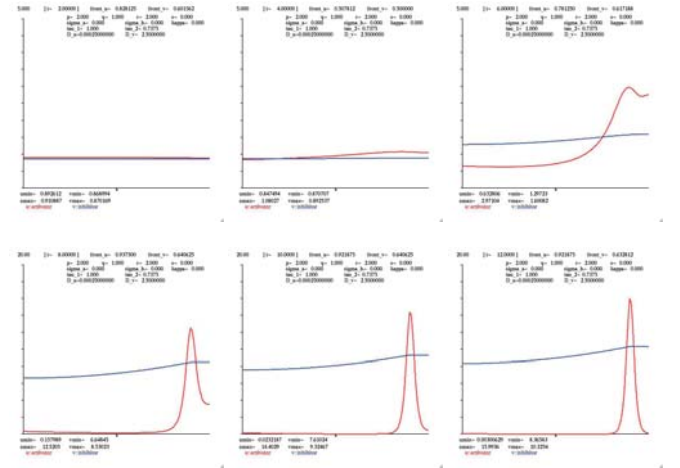
3 Background

A. Gierer and H. Meinhardt proposed the following reaction-diffusion system comprised of an activator $A = A(x, t)$ and an inhibitor $H = H(x, t)$ as the morphogenetic model of organisms:

$$\begin{cases} \frac{\partial A}{\partial t} = \varepsilon^2 \Delta A - \mu_a(x)A + \rho_a(x) \frac{A^p}{H^q} + \sigma_0(x) & \text{in } \Omega, \\ \tau \frac{\partial H}{\partial t} = D \Delta H - \mu_h(x)H + \rho_h(x) \frac{A^r}{H^s} & \text{in } \Omega. \end{cases} \quad (\text{GM})$$

Change in cells begin at the place of higher activator concentration. Stationary solutions of the shadow system ($D = +\infty, \sigma_0 \equiv 0$) are given by $(A(x), H(x)) \equiv (\xi^{q/(p-1)} u(x), \xi)$, where u solves

$$\varepsilon^2 \Delta u - \mu_a(x)u + \rho_a(x)u^p = 0.$$



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Heat kernel estimates for Markov processes associated with perturbed Dirichlet forms

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1 Notations and definitions

- $(\mathcal{E}, \mathcal{D}(\mathcal{E}))$: Dirichlet form on $L^2(\mathbb{R}^d)$ as follows:
 $\mathcal{E}(u, u) = \iint_{\mathbb{R}^d \times \mathbb{R}^d} (u(y) - u(x))^2 J(x, y) dx dy$,
 $\mathcal{D}(\mathcal{E}) = \{u \in L^2(\mathbb{R}^d) : \mathcal{E}(u, u) < \infty\}$
 $\frac{c_1}{|x-y|^d \phi(|x-y|)} \leq J(x, y) = J(y, x) \leq \frac{c_2}{|x-y|^d \phi(|x-y|)}$ for some positive constants c_1 and c_2 .
 $\phi(r) = r^\alpha$ or $r^\alpha \exp(mr)$ ($0 < \alpha < 2$, $m > 0$)
- $\mathcal{D}_e(\mathcal{E})$: extended Dirichlet space of $\mathcal{D}(\mathcal{E})$
- $\{P_t\}_{t \geq 0}$: semigroup associated with $(\mathcal{E}, \mathcal{D}(\mathcal{E}))$
- $\{X_t\}_{t \geq 0}$: Hunt process associated with $(\mathcal{E}, \mathcal{D}(\mathcal{E}))$
- $p(t, x, y)$: transition density function of $\{X_t\}$ or equivalently, the integral kernel of $\{P_t\}$.
- $G(x, y) = \int_0^\infty p(t, x, y) dt$.
- μ : positive Radon smooth measure on \mathbb{R}^d
- A_t^μ : positive continuous additive functional corresponding to μ .
- \mathcal{K} : the set of Kato class measure

Assumption 1. Dirichlet form $(\mathcal{E}, \mathcal{D}(\mathcal{E}))$ is transient, that is, $G(x, y) < \infty$ for $x \neq y$.

Definition 2. Let μ be the positive Radon smooth measure. The measure μ belongs to \mathcal{K}_∞ if $\mu \in \mathcal{K}$ and for arbitrary $\varepsilon > 0$, there exist positive constant $\delta > 0$ and compact set K such that

$$\sup_{x \in \mathbb{R}^d} \int_{K^c \cup B} G(x, y) \mu(dy) < \varepsilon$$

where B is an arbitrary set that satisfies $B \subset K$ and $\mu(B) < \delta$.

In the sequel we consider $\mu \in \mathcal{K}_\infty$.

Assumption 3. (i) The extended Dirichlet space $\mathcal{D}_e(\mathcal{E})$ is compactly embedded into $L^2(\mu)$.

$$(ii) \int_{\mathbb{R}^d \times \mathbb{R}^d} G(x, y) \mu(dx) \mu(dy) < \infty.$$

Remark 4. From [8] it is known that (i) of Assumption 3 holds when $\phi(r) = r^\alpha$.

2 Main result

Chen and the coauthors gave the two sided estimates of the transition density function $p(t, x, y)$ in [2, 3]. Here, we consider the Schrödinger form as follows:

$$\mathcal{E}^\mu(u, u) := \mathcal{E}(u, u) - \int_{\mathbb{R}^d} u^2 d\mu$$

Denote the semigroup associated with \mathcal{E}^μ by $\{P_t^\mu\}$. Then, it is known that $\{P_t^\mu\}$ admits the integral kernel $p^\mu(t, x, y)$ defined

on $(0, \infty) \times \mathbb{R}^d \times \mathbb{R}^d$ by [1]. We are going to consider the sufficient condition for $p^\mu(t, x, y)$ to have the same estimates as $p(t, x, y)$ does.

Definition 5. The measure μ is called gaugeable if it holds that $\sup_{x \in \mathbb{R}^d} \mathbb{E}^x[\exp(A_\infty^\mu)] < \infty$.

Set $h(x) = \mathbb{E}^x[\exp(A_\infty^\mu)]$. Since h is P_t^μ -excessive function, we can define h -transformed semigroup $\{P_t^{\mu, h}\}$ on $L^2(h^2 dx)$. Moreover, following the arguments of [4, 7], we can describe the corresponding Dirichlet form $(\mathcal{E}^{\mu, h}, \mathcal{D}(\mathcal{E}^{\mu, h}))$ as follows:

$$\mathcal{E}^{\mu, h}(u, u) = \iint_{\mathbb{R}^d \times \mathbb{R}^d} (u(y) - u(x))^2 J(x, y) h(x) h(y) dx dy$$

$$\mathcal{D}(\mathcal{E}^{\mu, h}) = \mathcal{D}(\mathcal{E})$$

Noting that $1 \leq h(x) \leq c_3$ for some positive constant c_3 , we can obtain the main result.

Theorem 6. Suppose $\mu \in \mathcal{K}_\infty$ be gaugeable. Under Assumptions 1 and 3, $p^\mu(t, x, y)$ has the same two-sided estimates as $p(t, x, y)$ does.

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On the enhancements to the Milnor numbers of a class of mixed polynomials

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The enhancement to the Milnor number $\lambda(K)$ is an invariant of a fibered link K on the 3-sphere S^3 defined via a vector field $\xi(K)$ on S^3 [1].

Construction of $\xi(K)$

- On $N(K)$, it is homotopic to $r(\frac{\partial}{\partial \theta}) + (1-r^2)\frac{\partial}{\partial \phi}$, where $N(K)$ is a tubular neighborhood of K with coordinates (r, θ, ϕ) .
- On $E(K) := S^3 \setminus \text{Int}N(K)$, it is a transverse field to the fiber surfaces of the fibration.
- On K , it is the tangent field of K .

Definition of $\lambda(K)$

$$\lambda(K) := \text{link}(\Delta^+(K), \Delta^-(K)) \in \mathbb{Z}$$

where $\Delta^\pm(K) = \{x \in S^3 \mid \xi(K)(x) = \pm t\psi(x), t > 0\}$ and ψ is a vector field which is homotopic to the tangent vector field of the Hopf fibration.

Remark.

We assume that $\xi(K)$ and ψ are vector fields in general position, so $\Delta^\pm(K)$ are 1-manifolds in S^3 .

In this poster, we study the following type of mixed polynomials

$$f(z, \bar{z}) := \prod_{j=1}^{m_+} (z_1^p + \alpha_j z_2^q) \prod_{j=m_++1}^{m_++m_-} \overline{(z_1^p + \alpha_j z_2^q)},$$

where $m_+ > m_-$, $\alpha_j \neq \alpha_{j'}$ ($j \neq j'$) and $\overline{z_1^p + \alpha_j z_2^q}$ represents the complex conjugate of $z_1^p + \alpha_j z_2^q$.

$f(z, \bar{z})$ satisfies the strong non-degeneracy condition. M. Oka showed the existence of the Milnor's fibration under this condition [2]. Remark that $f(z, \bar{z})$ is a special case of forms $(f\bar{g}, O)$ studied by A. Pichon and J. Seade [3, 4].

Thus $K_f := S_\varepsilon^3 \cap f^{-1}(0)$ is an oriented fibered link in S_ε^3 centered at the origin $O \in \mathbb{C}^2$ of radius ε .

The main theorem is the following:

Main Theorem For any $k \in \mathbb{Z}$, there exists a mixed polynomial $f(z, \bar{z})$ whose Milnor fibration $f/|f|: S_\varepsilon^3 \setminus K_f \rightarrow S^1$ satisfies $\lambda(K_f) = k$.

The sketch of the proof

The proof of Main Theorem consists of the following two steps.

[1] Construction of $\xi(K_f)$

The fibration is produced by the S^1 -action on S_ε^3 given by $(z_1, z_2) \mapsto (s^q z_1, s^p z_2)$ with parameter $s \in S^1$.

Let B be the orbit space of S_ε^3 under the S^1 -action. To construct $\xi(K_f)$, we prepare a vector field $\eta(K_f)$ on B as follows.

- The North and South poles are repeller points.
- There exist $m_+ + m_-$ mutually disjoint disks on B and, on each disk, $\eta(K_f)$ has exactly one attractor point and one saddle point.
- $\eta(K_f)$ has no zero outside these disks, N and S .

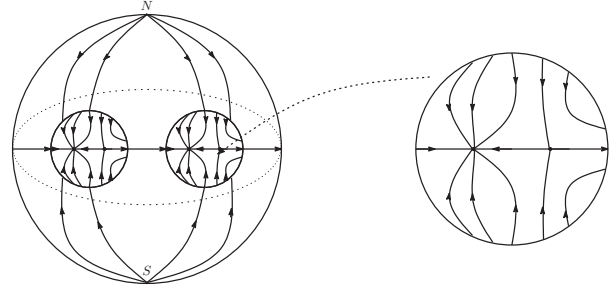


Figure 1: The vector field $\eta(K_f)$ and a 2-disk which has one attractor and one saddle point

Now we construct $\xi(K_f)$ as follows. On $N(K_f)$, set $\xi(K_f)$ such that it satisfies the following properties:

- The image of $N(K_f)$ by the orbit map is a union of disjoint 2-disks $D_1^2, \dots, D_{m_++m_-}^2$ each of which has only one attractor point.
- The image of $\xi(K_f)$ by the orbit map is $\eta(K_f)$ on $D_1^2, \dots, D_{m_++m_-}^2$.

On $E(K_f)$, $\xi(K_f)$ is constructed by a perturbation of ψ such that it is transverse to the fiber surfaces and its image by the orbit map is $\eta(K_f)$ on $B \setminus \bigcup_{j=1}^{m_++m_-} D_j^2$.

[2] Calculation of $\lambda(K_f)$

By the construction of $\xi(K_f)$, $\Delta^+(K_f)$ and $\Delta^-(K_f)$ are the preimages of zero points of $\eta(K_f)$.

We set

$\Delta_N^+, \Delta_S^+, \Delta_{\text{saddle}}^+$: the link components of $\Delta^+(K_f)$ corresponding to N, S and the saddle points.

K^+, K^- : the link corresponding to the holomorphic (resp. anti-holomorphic) factors of $f(z, \bar{z})$.

Then $\Delta^+(K_f) = \Delta_N^+ \cup \Delta_S^+ \cup \Delta_{\text{saddle}}^+ \cup K^+$ and $\Delta^-(K_f) = -K^-$. Thus $\lambda(K_f)$ can be calculated as follows:

$$\begin{aligned} \lambda(K_f) &= \text{link}(\Delta_N^+ \cup \Delta_S^+ \cup \Delta_{\text{saddle}}^+ \cup K^+, -K^-) \\ &= pm_- + qm_- - pq(m_+ + m_-)m_- + pqm_+m_- \\ &= (-pqm_- + p + q)m_-. \end{aligned}$$

We can check that $\lambda(K_f)$ realize all integers less than 2. The remaining integers can be realized by the mirror images of K_f .

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On the formal group of the Jacobian

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Formal Group

Let R be a commutative ring with the identity.

Let $X = (x_1, x_2, \dots, x_g)$ and $R[[X]]$ be the ring of formal power series with coefficients in R .

$$F(X, Y) = (F_1(X, Y), \dots, F_g(X, Y)), F_i(X, Y) \in R[[X, Y]]$$

is a formal group (or a group law) over R if it satisfies (1) and (2).

- (1) $F(Z, 0) = F(0, Z) = Z$
- (2) $F(F(X, Y), Z) = F(X, F(Y, Z))$

Examples of (a one-parameter) Formal Group

The additive formal group:

$$\mathbb{G}_a(x, y) = x + y$$

The multiplicative formal group:

$$\mathbb{G}_m(x, y) = (x + 1)(y + 1) - 1$$

For $f(x) \in R[[x]]$ such that $f(x) = ex + \dots$ (e is a unit in R), $F(x, y) = f^{-1}(f(x) + f(y)) \in R[[x, y]]$ is a formal group over R .

Elliptic Curve

An elliptic curve E over a field K is a nonsingular projective cubic curve, defined over K , with a K -rational point. $P = (x_1, \dots, x_n)$ is a K -rational point if $\forall i, x_i \in K$. We write $P \in E(K)$ if $E(P) = 0$. An elliptic curve over \mathbb{Q} can be expressed in Weierstrass form \hat{E} :

$$\hat{E}: y^2 + a_1xy + a_3y = x^3 + a_2x^2 + a_4x + a_6 \quad (a_i \in \mathbb{Q})$$

and by taking the minimal model of \hat{E} , a_i will be in \mathbb{Z} .

In general, a line through E meets E at 3 points. By using this fact, we can define an addition on an elliptic curve and the points on E form a group. This addition can be expressed as a formal group.

Honda's Theorem(1968)

Let E be an elliptic curve over \mathbb{Q} , E_w be the Weierstrass minimal model of E , $E_p = E_w \pmod{p}$ for a prime number p , and $E(x, y)$ be the formal group(A) given by the group law of E . For each prime number p , L_p is defined as:

$$L_p = \begin{cases} (1 - a_p p^{-s} + p^{1-2s})^{-1} & \text{if } E_p \text{ is of genus one;} \\ (1 - \varepsilon_p p^{-s})^{-1} & \text{if } E_p \text{ has an ordinary double point;} \\ 1 & \text{if } E_p \text{ has a cusp.} \end{cases}$$

where $a_p = p + 1 - \#E(\mathbb{F}_p)$, and $\varepsilon_p = 1$ if the tangents at the double point are rational over \mathbb{F}_p and $\varepsilon_p = -1$ if not. $a_n, l(x), L(x, y)$ are defined as:

$$\sum_{n=1}^{\infty} a_n n^{-s} = \prod_{p \in S} L_p(s)$$

$$l(x) = \sum_{n=1}^{\infty} n^{-1} a_n x^n$$

$$L(x, y) = l^{-1}(l(x) + l(y))$$

Let S be any set of prime numbers which does not contain 2 (resp. 3) if E_2 (resp. E_3) has genus one with $a_2 = \pm 2$ (resp. $a_3 = \pm 3$), and $\mathbb{Z}_S = \bigcap_{p \in S} (\mathbb{Z}_p \cap \mathbb{Q})$. Then

$L(x, y)$ is a formal group(B) over \mathbb{Z}

$$L(x, y) \cong E(x, y) \quad \text{over } \mathbb{Z}_S$$

*The restriction about S is removed by his paper in 1970.

In this theorem, there are two ways of the construction of a formal group; (A) and (B).

Hyperelliptic Curve

Let K be a field, $g \geq 1$ be an integer and let $h(x)$ and $f(x)$ be polynomials with coefficients in K such that $\deg f = 2g + 1$ and $\deg h \leq g$. Assume f is monic. The curve C given by the equation

$$C: y^2 + h(x)y = f(x)$$

is a hyperelliptic curve of genus g if it is nonsingular for all $x, y \in K$.

If $g = 1$, C is an elliptic curve. So the hyperelliptic curve is a generalisation of the elliptic curve.

If $g \geq 2$, a line meets C at more than 3 points, so we cannot define the addition on C like an elliptic curve. Instead, we use the divisor.

Jacobian Variety

Let C be a hyperelliptic curve defined over a field K . For each point $P \in C(K)$, define a formal symbol $[P]$. A divisor D on C is a finite linear combination of such symbols with integer coefficients:

$$D = \sum_j a_j [P_j], \quad a_j \in \mathbb{Z}$$

Let f be a function on C . We define the degree of a divisor and the divisor of f by

$$\deg(\sum_j a_j [P_j]) = \sum_j a_j (\in \mathbb{Z})$$

$$\text{div}(f) = \sum_{P \in C(K)} \text{ord}_P(f) [P]$$

The divisor of a function on C is called a principal divisor. It is known that the degree of a principal divisor is always zero.

The group of divisors of degree 0 modulo principal divisors is called the Jacobian variety J of C .

Generalisation of Honda's Theorem (Freije (1993), Sairaiji (2010))

Let C be a complete nonsingular algebraic curve of genus g over a field K of $\text{ch}(K) = 0$. Assume P is a K -rational point on C which is not a Weierstrass point. Choose a basis η_1, \dots, η_g for the space of the holomorphic differentials of C such that the K -expansion of η_i with respect to the parameter t satisfies

$$\eta_i \equiv (-t)^{i-1} \pmod{t^g dt}$$

Let $l_i(t) = \int \eta_i$ satisfying $l_i(0) = 0$ ($1 \leq i \leq g$) and let $S(X) = (s_1(X), \dots, s_g(X))$ where $s_i(X)$ is the i th symmetric function on g letters. Define $L(X) = (L_1(X), \dots, L_g(X))$ by

$$L_i(S(T)) = l_i(t_1) + \dots + l_i(t_g)$$

Then the formal group of J, \hat{J} , is:

$$\hat{J} = L^{-1}(L(X) + L(Y))$$

Freije showed that the formal group of J can be constructed by using the expansion of a holomorphic differential of J by a local parameter at the zero. This is a generalisation of a formal group(B) of Honda's theorem. She also proved that $J_l(X, Y) \in \mathbb{Z}_p[[X, Y]]$ for all primes $p \notin S$, where S is some finite set of primes, in the case of the modular curve $X_0(l)$ (l a prime).

In his 2010 paper, Sairaiji constructed a formal group \hat{J} by expanding the addition formula of J by a system of local parameters at the zero of $J \times J$ in the case of the hyperelliptic curve and showed that the ring generated by the coefficients of \hat{J} is essentially generated over \mathbb{Z} by the coefficients of C . This is a generalisation of a formal group(A) of Honda's theorem.

On the density of some sequences of integers

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1 Notation and Definition of Density

Let $n = p_1^{\alpha_1} \cdots p_r^{\alpha_r}$ be the prime factorization of a positive integer n .

Def 1.1. Define the **excess** of n be $(\alpha_1 - 1) + \cdots + (\alpha_r - 1)$.

An integer with excess 0 is said to be **square-free**.

Let A denote \mathbb{Z} or $\mathbb{F}_q[x]$ for some prime power $q = p^e$ and K denote the fraction field of A . Define

$$\begin{aligned} \text{Box} &= \text{Box}(B_1, \dots, B_l) \\ &:= \{(a_1, \dots, a_l) \in A^l; |a_i| \leq B_i \text{ for all } i\}. \end{aligned}$$

For $E \subseteq A^l$, define that the density of E is the limit (if it exists)

$$\mu(E) := \lim_{B_1, \dots, B_l \rightarrow \infty} \frac{\#(E \cap \text{Box})}{\#\text{Box}}.$$

Also, define the density in weaker sense,

$$\mu_l(E) := \lim_{B_1, \dots, B_{n-1} \rightarrow \infty} \lim_{B_l \rightarrow \infty} \frac{\#(E \cap \text{Box})}{\#\text{Box}}.$$

This has the effect of considering only boxes in which the l th dimension is large relative to the others.

2 Density of $E_k \subset \mathbb{Z}$

Let E_k denote a set of positive integers with excess k .

If $A = \mathbb{Z}$, the density of $E_k \subset \mathbb{Z}$ is the limit (if it exists)

$$\mu(E_k) = \lim_{n \rightarrow \infty} \frac{\#(E_k \cap \{1, 2, \dots, n\})}{n},$$

which is the limiting "probability" that an integer from 1 to n is in E_k .

Thm 2.1. (Rényi) *The set E_k has a density d_k and that the generating function of the sequence $\{d_k\}$ is given by*

$$\sum_{k=0}^{\infty} d_k z^k = \prod_{p:\text{prime}} \left(1 - \frac{1}{p}\right) \left(1 + \frac{1}{p-z}\right). \quad (1)$$

Substituting $z = 0$ into (1), we obtain the special case

$$d_0 = \prod_p \left(1 - \frac{1}{p^2}\right) = \frac{1}{\zeta(2)} = \frac{6}{\pi^2}, \quad (2)$$

which is well-known result that the density of square-free integers.

Substituting $z = 1$ into (1), one sees that $\sum_k d_k = 1$, so that d_k can be considered as elements of a probability distribution.

3 The analogue of Rényi's result

Let $f = \pi_1^{\alpha_1} \cdots \pi_l^{\alpha_l}$ be a monic polynomial with prime factorization, and define the excess of f to be $(\alpha_1 - 1) + \cdots + (\alpha_r - 1)$, just as integers. Let $e_{l,k}$ be the number of monic polynomials of degree l and excess k . Define $d_{l,k} = e_{l,k}/q^l$, which is the probability that a monic polynomial of degree l has excess k . Then define the analogue of the density to be the limiting "probability" as

$$d_k = \lim_{l \rightarrow \infty} d_{l,k}. \quad (3)$$

Let $Nf = q^{\deg f}$ be the norm of f , which is the cardinality of the residue ring $\mathbb{F}_q[x]/(f)$.

Thm 3.1. (Morrison) *The generating function $D(z)$ of the sequence $\{d_k\}$ has a factorization over the prime polynomials given by*

$$\begin{aligned} D(z) &= \sum_{k=0}^{\infty} d_k z^k \\ &= \prod_{\pi:\text{prime}} \left(1 - \frac{1}{N\pi}\right) \left(1 + \frac{1}{N\pi - z}\right). \end{aligned} \quad (4)$$

4 Square-free Values over A

Now suppose $f(x) \in \mathbb{Z}[x]$ is a polynomial and let $E = E_0$ be the set of $n \in \mathbb{Z}_{\geq 0}$ for which $f(n)$ is square-free. Then one guesses that

$$\mu(E) = \prod_{p:\text{prime}} \left(1 - \frac{c_p}{p^2}\right), \text{ where } c_p := \#\{n \in [0, p^2 - 1]; p^2 | f(n)\}$$

When $\deg f \leq 3$, this guess is correct. For general f with $\deg f \geq 4$, it's unknown whether this is correct. But Granville and Poonen showed the following theorems.

Thm 4.1. (Square-free values over \mathbb{Z}) *Assume the abc-conjecture. Let $f \in \mathbb{Z}[x_1, \dots, x_l]$ be a polynomial that square-free as an element of $\mathbb{Q}[x_1, \dots, x_l]$ and suppose that x_l appears in f . For each prime p , let*

$$E_f := \{a \in \mathbb{Z}^l; f(a) \text{ is square-free}\},$$

$$c_p := \#\{x \in (\mathbb{Z}/p^2)^l; f(x) = 0 \text{ in } \mathbb{Z}/p^2\}.$$

Then the density of E_f is

$$\mu(E_f) = \prod_p \left(1 - \frac{c_p}{p^{2l}}\right). \quad (5)$$

Remark This time the *abc*-conjecture is used to bound the number of polynomial values divisible by the square of a large prime.

Thm 4.2. (Square-free values over $\mathbb{F}_q[t]$) *Let $A = \mathbb{F}_q[t]$. Let $f \in A[x_1, \dots, x_l]$ be a polynomial that square-free as an element of $K[x_1, \dots, x_l]$. For each non-zero prime $\mathfrak{p} \subset A$, let*

$$E_f := \{a \in A^l; f(a) \text{ is square-free}\},$$

$$c_{\mathfrak{p}} := \#\{x \in (A/\mathfrak{p}^2)^l; f(x) = 0 \text{ in } A/\mathfrak{p}^2\}.$$

Then the density of E_f is

$$\mu(E_f) = \prod_{\mathfrak{p}} \left(1 - \frac{c_{\mathfrak{p}}}{|\mathfrak{p}|^{2l}}\right). \quad (6)$$

5 Square-free Values over function fields

Thm 5.1. (Square-free Values over other rings of functions) *Let S be a finite nonempty set of closed points of smooth, projective, geometrically integral curve X over \mathbb{F}_q . Define the affine curve $U := X \setminus S$, and let A be the ring of regular functions on U . Then there is the analogous result of Thm4.2.*

The proof of Thm5.1 is the analogue of the proof of Thm4.2.

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Homogeneous Reinhardt domains of Stein in the complex n-space

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Algebraic isomorphisms

An analytic automorphism $(z_i) \mapsto (w_i)$ of $(\mathbb{C}^*)^n$ is called an algebraic automorphism, if whose components are given by Laurent monomials, that is, of the form

$$w_i = \alpha_i z_i^{a_{i1}} \cdots z_n^{a_{in}} \quad (1 \leq i \leq n)$$

where $(a_{ij}) \in GL(n, \mathbb{Z})$ and $(\alpha_i) \in (\mathbb{C}^*)^n$.

Suppose D_1 and D_2 are domains in \mathbb{C}^n and an analytic isomorphism $\varphi : D_1 \mapsto D_2$ is induced by an algebraic automorphism. Then the isomorphism φ is said to be an algebraic isomorphism, and two domains D_1, D_2 are called algebraically equivalent.

Bounded Reinhardt domains

A domain in \mathbb{C}^n is called a Reinhardt domain, if it is stable under rotations around the coordinate axis. As is well known, Poincaré showed that there is no analytic isomorphism of the polydisc to the unitball in \mathbb{C}^2 . By the way, both of them are Reinhardt domains. This is sharp contrast to the Riemann mapping theorem in \mathbb{C} , i.e. analytic isomorphisms of several complex variables are considered to be high rigid. Sunada and Shimizu generalized by far over the Poincaré example:

Theorem 1. (Shimizu) Two bounded Reinhardt domains in \mathbb{C}^n are holomorphically equivalent, if and only if they are algebraically equivalent.

Homogeneous Reinhardt domains

Unfortunately it is difficult to generalize Theorem 1 in unbounded Reinhardt domains, because in such a case we cannot use the classical Cartan theorem. And so, extending algebraically equivalence to unbounded cases, we want to classify homogeneous Reinhardt domains of Stein. This classification is conjectured to be as follows:

Conjecture 1. For a homogeneous Reinhardt domain D to be Stein in \mathbb{C}^n , there exist integers n_1, \dots, n_k and non-negative integers l, m with $n_1 + \dots + n_k + l + m = n$ such that

$$D \stackrel{alg}{\cong} B_{n_1} \times \cdots \times B_{n_k} \times \mathbb{C}^l \times (\mathbb{C}^*)^m,$$

where B_{n_i} is the complex n_i -unit ball.

Shimizu proved this conjecture in the case of which D is bounded, but in the unbounded case it is still remains. The next result in my master's thesis is a partial solution of the problem in the latter case. I announced it in Several Complex Variables Winter Seminar 2011 at Hiroshima.

Theorem 2. If a Reinhardt domain D in $(\mathbb{C}^*)^n$ is homogeneous and Stein, then the domain D is identified with $(\mathbb{C}^*)^n$.

Future study

Now I shall research on Conjecture 1 under the opposite condition to it of Theorem 2, that is, assume a domain contains the origin. In addition, for the sake of simplicity, let the complex dimension be equal to 3:

Conjecture 2. For a homogeneous Reinhardt domain D in \mathbb{C}^3 that is Stein and contains the origin, there exist non-negative integers n_1, n_2, n_3 and m with $n_1 + n_2 + n_3 + m = 3$ such that

$$D \stackrel{alg}{\cong} B_{n_1} \times B_{n_2} \times B_{n_3} \times \mathbb{C}^m$$

where B_{n_i} is the complex n_i -unit ball.

I think the notion of Liouville foliation, which was introduced by Shimizu in order to analyze tube domains, would play a key role in a solution.

Definition 1. Let M be a complex manifold. A collection $\{\Sigma_\alpha\}_{\alpha \in A}$ of subsets of M is called a Liouville foliation on M if the following conditions are satisfied:

- (L1) If $\alpha_1 \neq \alpha_2$, then $\Sigma_{\alpha_1} \cap \Sigma_{\alpha_2} = \emptyset$;
- (L2) $\bigcup_{\alpha \in A} \Sigma_\alpha = M$;
- (L3) For each Σ_α , any bounded holomorphic function on M takes a constant value on Σ_α ;
- (L4) $\alpha_1, \alpha_2 \in A$ with $\alpha_1 \neq \alpha_2$, there exists a bounded holomorphic function h on M such that the constant values of h on Σ_{α_1} and Σ_{α_2} are different.

In the above definition, if we exchange 'holomorphic' for 'plurisubharmonic', then a Liouville foliation is said to be a plurisubharmonic Liouville foliation. Furthermore, we can fuse these two notions for the foliations on M : Suppose a mapping $h = (f_1, \dots, f_p, g_1, \dots, g_q) : M \rightarrow \mathbb{C}^p \times \mathbb{R}^q$ satisfies following conditions:

- (1) $f_i : M \rightarrow \mathbb{C}$ is a bounded holomorphic function,
- (2) $g_j : M \rightarrow \mathbb{R}$ is a bounded plurisubharmonic function.

Then h is called a Liouville mapping of type (p, q) . By this, we can think of Liouville foliation with type (p, q) in the same way.

Calibrated Submanifolds

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SYZ Conjecture

Mirror symmetry is a mysterious relationship between pairs of Calabi-Yau 3-folds X, \hat{X} . It was discovered by physicists working in string theory.

SYZ Conjecture explains mirror symmetry in terms of dual fibrations $f: X \rightarrow B, \hat{f}: \hat{X} \rightarrow B$ with special Lagrangian fibers.

Similar statement also holds for pairs of G_2 manifolds in terms of dual coassociative fibrations.

Calibrated Geometry

Definition .Let (M, g) be an m -dim Riemannian manifold and φ be the closed k -form on M ($1 \leq k \leq m$).

We say φ is **calibration** on M if for any $p \in M$ and any oriented k -dim subspace of $T_p M$ we have

$$\varphi|_V \leq \text{vol}_V.$$

Let N be the k -dim oriented submanifold of M . We say N is a **calibrated submanifold**(φ -submanifold) of M if

$$\varphi|_N = \text{vol}_N.$$

Examples

$\text{Hol}(g) (\subset)$	$U(m)$	$SU(m)$	G_2
(M, g)	Kähler	Calabi-Yau	G_2
φ	$\omega^k/k!$ (ω :Kähler form)	$\text{Re}(\Omega)$ (Ω :hol. vol. form)	$\varphi \in \Omega^3$ $*\varphi \in \Omega^4$ (φ : G_2 -structure)
φ -sub manifolds	k -dim complex submfds	special Lagrangian submfds	(co)associative submfds

Construction of Calibrated Submfds

We focus on that of SL submfds.

Fact .Let (M, J, ω, Ω) be an m -dim CY mfd and L be a real m -dim oriented submfd of M . Then,

L is a SL submfd of $M \Leftrightarrow \omega|_L = 0$ (i.e. L : Lagrangian), $\text{Im}\Omega|_L = 0$.

From this fact, we know SL submfds are **characterized by vanishing of forms**. This fact is very useful for the construction.

◆Graph in \mathbb{C}^m

For any $f \in C^\infty(\mathbb{R}^m)$, $\text{graph}(df) \subset T^*\mathbb{R}^m = \mathbb{C}^m$ is a Lagrangian submfd.

Solving a second order elliptic PDE, we get SL submfds.

◆Evolution Equations in \mathbb{C}^m

$\text{Re}(\Omega)$ is "maximal" on SL submfds.

Let $P \subset \mathbb{C}^m$ be an $(m-1)$ -dim submfd.

Solving ODE for $\{\phi_t: P \rightarrow \mathbb{C}^m\}_{t \in \mathbb{R}}$ to flow $\phi_t(P)$ in the direction in which $\text{Re}(\Omega)$ is "largest", we get SL submfds.

◆SL submanifolds with Large symmetry

Assume that a compact Lie group G acts on a CY mfd M preserving CY structure and there exists a **moment map** μ for the G -action.

Connected G -invariant Lagrangian submfds must be in the level set of μ .

Find G -inv. SL submfds in the level set of μ .

When $G = T^{m-1}$ and M is simply connected, there exists a T^{m-1} -inv. function $f \in C^\infty(M)$ and

$$(\mu, f): M \rightarrow U \overset{\text{open}}{\subset} (\mathfrak{t}^m)^* \times \mathbb{R}$$

is a **SL fibration**.

Remark .Similar methods can be considered for coassociative submfds from the following fact.

Fact .Let (Y, φ) be a torsion-free G_2 mfd and L be a 4-dim oriented submfd of Y . Then, L is a coassociative submfd of $M \Leftrightarrow \varphi|_L = 0$.

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On Legendrian minimal submanifolds in Sasakian manifolds

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Introduction

In 1990's, Y.-G.Oh introduced the notion of *Hamiltonian-minimal (H-minimal)* Lagrangian submanifolds in Kähler manifolds ([2]). This is a nice extension of the notion of *minimal* submanifold, and has been studied by many researchers. One of the main concern is H-stability of an H-minimal Lagrangian submanifold in a specific Kähler manifold.

On the other hand, there is a notion of *Sasakian manifolds* which is an odd-dimensional counterpart of Kähler manifolds. In Sasakian manifolds, we consider *Legendrian-minimal (L-minimal)* Legendrian submanifolds which corresponds to H-minimal Lagrangian manifolds in Kähler manifolds. We call L-minimal Legendrian submanifolds *Legendrian-stable (L-stable)* if the second variation is non-negative for any Legendrian deformations. In this presentation, we mention some results for the construction and the L-stability of L-minimal Legendrian submanifolds in specific Sasakian manifolds.

Definition

An odd-dimensional Riemannian manifold (M^{2n+1}, g) is called *Sasakian manifold* if its metric cone $(M \times \mathbf{R}_{>0}, r^2g + dr^2)$ admits a Kähler structure. A Sasakian manifold is a contact manifold with the contact structure (ϕ, ξ, η, g) , where η is a contact 1-form, ξ is the characteristic vector field, ϕ is an $(1, 1)$ -tensor.

Definition 1. Let $(M^{2n+1}, \phi, \xi, \eta, g)$ be a Sasakian manifold, and L^n a n -dimensional manifold.

(1) An immersion $\iota : L^n \rightarrow M^{2n+1}$ is called *Legendrian* : $\iff \iota^*\eta = 0$.

(2) An Legendrian immersion ι is called *L-minimal* : $\iff \operatorname{div} \phi H = 0$, where H is the mean curvature vector of ι .

Remark 1. A deformation $\{\iota_t\}$ of ι is called *Legendrian deformation* if it preserves the Legendrian condition, i.e., $\iota_t^*\eta = 0$. The L-minimal Legendrian immersion is a critical point of the volume functional under the Legendrian deformation. Thus the L-minimal immersion is an extension of minimal immersion.

A compact L-minimal immersion ι is called *Legendrian stable (L-stable)* if the second variation is non-negative for any Legendrian deformations.

L-minimal immersion into $\mathbf{R}^{2n+1}(-3)$

The odd-dimensional Euclidean space admits the standard Sasakian structure. We denote this $\mathbf{R}^{2n+1}(-3)$.

Theorem 1 ([1]). Let $\iota : L^n \rightarrow \mathbf{R}^{2n+1}(-3)$ be a Legendrian immersion. Assume that L^n lies in a cylinder

$$N_{\mathbf{x}_0}^{2n}(r) := \{x \in \mathbf{R}^{2n+1} | g(x - \mathbf{x}_0, x - \mathbf{x}_0) - \eta(x - \mathbf{x}_0)^2 = r^2\},$$

where $\mathbf{x}_0 := (x_0^1, \dots, x_0^n, y_0^1, \dots, y_0^n, z_0)$ is a constant vector in \mathbf{R}^{2n+1} and r is a constant. If ι has the parallel mean curvature in $N_{\mathbf{x}_0}^{2n}$, then ι is L-minimal in $\mathbf{R}^{2n+1}(-3)$.

In the case of $n = 1$, L-minimal Legendrian curve in $\mathbf{R}^3(-3)$ is the only the geodesic or the curve which is minimal in the cylinder. In general, an immersion into $\mathbf{R}^{2n+1}(-3)$ which lies in some cylinders and minimal in the cylinder is called *1-type*. By Theorem 1, we have

Corollary 1. An 1-type Legendrian immersion is L-minimal.

Remark 2. Similar result of Theorem 1 for H-minimal Lagrangian submanifolds in \mathbf{C}^n was obtained by Y.-G.Oh in [2].

L-minimal curves in 3-dim Sasakian space forms

A Sasakian manifold is called *Sasakian space form* if it has constant ϕ -sectional curvature ($= c$). We denote Sasakian space form $M^{2n+1}(c)$.

Theorem 2 ([1]). Let γ be a compact L-minimal Legendrian curve in a 3-dimensional Sasakian space form $M^3(c)$. Then γ is L-stable if and only if it satisfies $\lambda_1 \geq c + 3 + h^2$, where λ_1 is the first eigenvalue of the Laplace-Beltrami operator Δ acting on $C^\infty(\gamma)$, and h is the curvature of γ .

Example 1 ([1]). (1) The 3-dimensional unit sphere $S^3(1)$: All of closed L-minimal Legendrian curves in $S^3(1)$ are obtained by Schoen-Wolfson and

Iriyeh:

$$\gamma(s) = \frac{1}{\sqrt{p+q}} \left(\sqrt{q} e^{\sqrt{-1}\sqrt{\frac{q}{p}}s}, \sqrt{-1}\sqrt{p} e^{-\sqrt{-1}\sqrt{\frac{p}{q}}s} \right), \quad s \in [0, 2\pi\sqrt{pq}],$$

where (p, q) is a pair of relatively prime positive number. They are torus knots of type (p, q) , and in the case of $(p, q) = (1, 1)$, it is minimal. By Theorem 2, these curves are all L-unstable.

(2) $\mathbf{R}^3(-3)$: An L-minimal Legendrian curve in $\mathbf{R}^3(-3)$ is the only the geodesic ($h = 0$) or the following:

$$\gamma(s) = \mathbf{x}_0 + \left(\frac{2}{h} \cos hs, \frac{2}{h} \sin hs, -\frac{2}{h}s + \frac{1}{h^2} \sin 2hs + \frac{2y_0}{h} \cos hs \right),$$

where $\mathbf{x}_0 = (x_0, y_0, z_0)$ and s is the arc-length parameter. These curves are *helices*, and so *not periodic*. If we cut the curve with the length l , by Theorem 2, L-minimal curve in $\mathbf{R}^3(-3)$ is L-stable if and only if $0 \leq h \leq \pi/l$.

(3) The special linear group $SL(2, \mathbf{R}) = M^3(-7)$: We choose a global coordinate $(x, y, \theta) \in \mathbf{R} \times \mathbf{R}^+ \times S^1$ of $SL(2, \mathbf{R})$ by the Iwasawa decomposition. Then L-minimal curves in $SL(2, \mathbf{R})$ are given by:

$$\gamma(s) = \left(r \sin \mu(s) + x_0, r \left(\frac{h}{2} - \cos \mu(s) \right), \frac{\mu(s)}{2} - \frac{h}{2}s + \theta_0 \right),$$

where $r \in \mathbf{R}^+$ is a positive constant, h is the curvature of γ , and the smooth function $\mu : \mathbf{R} \rightarrow \mathbf{R}$ is a solution of the ODE $\dot{\mu}(s) = h - 2 \cos \mu(s)$. The above curve is closed if and only if $h = 2/\sqrt{1 - (m/k)^2}$ for some relatively positive integers $m, k \in \mathbf{Z}^+$ with $m/k < 1$. By Theorem 2, the closed curve is L-minimal if and only if $h = 2/\sqrt{1 - (1/k)^2}$. Thus there exist an closed L-stable L-minimal curves in $SL(2, \mathbf{R})$.

Some relations between H-minimal Lagrangian submanifold and L-minimal Legendrian submanifold

The notion of L-minimal Legendrian submanifold is closely related to the notion of H-minimal Lagrangian submanifold. The following facts are known.

Fact 1. Let $\pi : S^{2n+1}(1) \rightarrow \mathbf{C}P^n$ be the Hopf fibration, and $\bar{\iota} : \bar{L}^n \rightarrow \mathbf{C}P^n$ a Lagrangian submanifold. Then we have a Legendrian immersion $\iota : L^n \rightarrow S^{2n+1}(1)$ as the lift of $\bar{\iota}$. Moreover, ι is L-minimal if and only if $\bar{\iota}$ is H-minimal.

Fact 2. Let $\iota : L^n \rightarrow S^{2n+1}(1)$ be a Legendrian submanifold. Then the cone $C(L) := L^n \times \mathbf{R}_{>0}$ of L^n is Lagrangian submanifold in $S^{2n+1} \times \mathbf{R}_{>0}$. Moreover L is L-minimal if and only if $C(L)$ is H-minimal.

Remark 3. We can generalize Fact 1 for the canonical fibration $\pi : M^{2n+1} \rightarrow M^{2n}$. For example, $\pi : \mathbf{R}^{2n+1}(-3) \rightarrow \mathbf{C}^n$. We can generalize Fact 2 for any Sasakian manifolds.

L-unstability Theorem for L-minimal immersion into the unit sphere

In $\mathbf{C}P^n$, it is known that there exist a compact H-stable H-minimal Lagrangian submanifold. For example, the Clifford torus and the parallel Lagrangian submanifold (which means Lagrangian submanifold with the parallel second fundamental form) are all H-stable. In contrast to this fact, we have the following theorem for L-minimal immersion into $S^{2n+1}(1)$:

Theorem 3. There are no compact L-stable L-minimal Legendrian immersion into the unit sphere $S^{2n+1}(1)$.

Remark 4. For the case of minimal Legendrian submanifolds in $S^{2n+1}(1)$, L-unstability theorem has been proven by H.Ono. The case of $n = 1$ follows from Theorem 2.

To prove Theorem 3, we find a Legendrian vector field along ι which gives the Legendrian deformation $\{\iota_t\}$ with $\frac{d^2}{dt^2} \operatorname{Vol}(\iota_t) < 0$.

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On elliptic surfaces related to Beilinson' s Tate conjecture

$K_2(U)_{Q_p}^{(2)}$ and $H_{\text{et}}^2(U_{\bar{k}}, Q_p(2))^{G_k}$ are very important objects over a variety U !!

We consider the case that U is an open subscheme of an elliptic surface $\pi: E \rightarrow C$. Assume that $P_1, \dots, P_s \in C$ give all fibers of split type of $\pi: E \rightarrow C$. Put $D_i = \pi^{-1}(P_i)$, $D = \sum_{i=1}^s D_i$ and $U = E \setminus D$.

$$\begin{array}{ccc}
 K_2(U)_{Q_p}^{(2)} & & \\
 \downarrow c & \searrow \partial_{Q_p} & \\
 H_{\text{et}}^2(U_{\bar{k}}, Q_p(2))^{G_k} & \xrightarrow{\partial_{\text{et}}} & \bigoplus_{i=1}^s \vec{Q}_p[D_i].
 \end{array}$$

The surjectivity of c is called "Beilinson' s Tate conjecture"

The surjectivity of c follows from the surjectivity of ∂

Considering the relation between them is a historical problem.

We present examples of an elliptic surface over a field k whose open subscheme U satisfies Beilinson' s Tate conjecture but the surjectivity of ∂_Q !!!

Main Theorem

Theorem Let F be a field and $\text{ch}(F)$ the characteristic of F . Set $k = F(S)$, where S is an indeterminate element. Consider the rational elliptic surfaces defined by the Weierstrass forms over $k(t)$

$$E_1: y^2 = x^3 + x^2 + tx + t(1 - S^2), \quad \text{if } \text{ch}(F) \neq 2, 3$$

$$E_2: y^2 = x^3 + (S^2 - 1)t^2x^2 + t^3(t - 1), \quad \text{if } \text{ch}(F) = 3$$

$$E_3: y^2 + xy = x^3 + S^2x^2 + t^2(S - t), \quad \text{if } \text{ch}(F) = 2.$$

Then the map $\partial_Q: K_2(U)_Q^{(2)} \rightarrow \bigoplus_{i=1}^s Q[D_i]$ is not surjective for the above rational elliptic surfaces.

The Law of the Iterated Logarithm (LIL) for G -Brownian Motion

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Motivated by volatility uncertainty problems in finance, S. Peng introduced the notion of G -Brownian motion. Intuitively, G -Brownian motion is a Brownian motion whose variance is uncertain. The purpose of this study is a quasi-sure analysis of paths of G -Brownian motion.

Notation

- $\Omega := C_0([0, \infty); \mathbb{R})$
- B : the canonical process of Ω

A nonlinear heat equation

$$\frac{\partial u}{\partial t} - \frac{1}{2} \sup_{\gamma \in [\sigma_0, \sigma_1]} \left\{ \gamma^2 \frac{\partial^2 u}{\partial x^2} \right\} = 0 \quad \text{in } (0, \infty) \times \mathbb{R},$$

is called G -heat equation. Here $[\sigma_0, \sigma_1]$ is a given subset of $[0, \infty)$.

Using the viscosity solutions of G -heat equation, Peng [4] constructed a sublinear expectation, called G -expectation, under which B becomes a **G -Brownian motion**; the given subset $[\sigma_0, \sigma_1]$ expresses the volatility uncertainty of G -Brownian motion.

Denis–Hu–Peng [1] proved that G -expectation can be represented as the supremum of the linear expectations with respect to martingale

measures. The precise claim is as follows:

Let

W : B.m. on a prob. sp. (Ω, \mathcal{F}, P)

\mathcal{A} : the set of all $[\sigma_0, \sigma_1]$ -valued prog. m'ble proc. on $[0, \infty)$

P_θ : the law of $\int_0^\cdot \theta_s dW_s, \quad \theta \in \mathcal{A}$

Theorem. (Denis–Hu–Peng [1])

$$\mathbb{E}[X] = \sup_{\theta \in \mathcal{A}} E_{P_\theta}[X]$$

Define a pair of capacities (V, v) by

$$V(A) := \sup_{\theta \in \mathcal{A}} P_\theta(A), \quad v(A) := \inf_{\theta \in \mathcal{A}} P_\theta(A) \quad \text{for } A \in \mathcal{B}(\Omega).$$

A proposition holds **quasi-surely** if it holds outside a set N with $v(N^c) = 1$.

\leadsto It becomes possible to investigate pathwise properties of G -Brownian motion.

LIL for G -Brownian motion (Hariya–O. [3])

$$\sigma_0 \leq \limsup_{t \rightarrow \infty} \frac{B_t}{\sqrt{2t \log \log t}} \leq \sigma_1 \quad \text{quasi-surely.}$$

(The idea of the proof in [3])

Time-change formula

\exists B.m. β on (Ω, \mathcal{F}, P) s.t.

$$\int_0^t \theta_s dW_s = \beta \left(\int_0^t \theta_s^2 ds \right), \quad \forall t \geq 0, \quad P\text{-a.s.}$$

Combining this with LIL for β , we have

$$P_\theta \left(\sigma_0 \leq \limsup_{t \rightarrow \infty} \frac{B_t}{\sqrt{2t \log \log t}} \leq \sigma_1 \right) = 1, \quad \forall \theta \in \mathcal{A}.$$

(Another proof)

LDP of Schilder's type for G -B.m.:

- upper bound of LDP for $\{V(\varepsilon B \in \cdot), \varepsilon > 0\}$ (Gao–Jiang [2])
- lower bound of LDP for $\{v(\varepsilon B \in \cdot), \varepsilon > 0\}$

\Downarrow

Invariance principle of G -B.m. for LIL (Wu–Chen [5])

\Downarrow

LIL for G -B.m.

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Modified wave operator for the 2d nonlinear Schrödinger system with mass resonance

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1 Introduction

We study the asymptotic behavior of the solution to the following nonlinear Schrödinger system:

$$(2\text{-NLS}) \quad \begin{cases} i\partial_t u_1 + \frac{1}{2m_1} \Delta u_1 = \bar{u}_1 u_2, & t \in \mathbb{R}, x \in \mathbb{R}^2, \\ i\partial_t u_2 + \frac{1}{2m_2} \Delta u_2 = u_1^2, & t \in \mathbb{R}, x \in \mathbb{R}^2. \end{cases}$$

where $u_1(t, x), u_2(t, x) : \mathbb{R} \times \mathbb{R}^2 \rightarrow \mathbb{C}$; unknown complex valued functions. m_1, m_2 ; positive constants.

Weighted Sobolev spaces Let $s, m \in \mathbb{R}$,

$$H^{s,m} := \left\{ \psi \in \mathcal{S}' ; \|(1 + |x|^2)^{m/2} (1 - \Delta)^{s/2} \psi\|_{L^2} < \infty \right\}.$$

2 Known results (Single case)

Nonlinear Schrödinger equation:

$$(NLS) \quad i\partial_t u + \frac{1}{2m} \Delta u = f(u), \quad t \in \mathbb{R}, x \in \mathbb{R}^2.$$

- $f(u) = |u|^{p-1}u$ $p > 2 : u(t) \sim e^{\frac{i}{2m}t\Delta}u_+$,
 $p \leq 2 : u(t) \not\sim e^{\frac{i}{2m}t\Delta}u_+$.
 Barab ('84), Y. Tsutsumi-Yajima ('84).
- $f(u) = |u|u$ $u(t) \sim e^{iS(t)}e^{\frac{i}{2m}t\Delta}u_+$,
 $S(t) := -m|\hat{u}_+(mx/t)| \log t$.
 Ozawa ('91), Ginibre-Ozawa ('93).
- $f(u) = u^2$ or \bar{u}^2 $u(t) \sim e^{\frac{i}{2m}t\Delta}u_+$.
 Moriyama-Tonegawa-Y. Tsutsumi ('03).
- $f(u) = |u|^2$ $u(t) \not\sim e^{\frac{i}{2m}t\Delta}u_+$.
 Shimomura ('05),
 Shimomura-Y. Tsutsumi ('06).

3 Results

Theorem 1. *Let $2m_1 = m_2$, $1 < \alpha < 2$. Then there exists $\varepsilon > 0$ with the following property:*

For any $u_{1+}, u_{2+} \in H^{0,\alpha}$ with $|\hat{u}_{1+}(\xi)| = \sqrt{2}|\hat{u}_{2+}(\xi)|$ ($\forall \xi \in \mathbb{R}^2$) and $\|\hat{u}_{1+}\|_{H^\alpha} + \|\hat{u}_{2+}\|_{H^\alpha} \leq \varepsilon$, (2-NLS) has a unique global solution $(u_1(t), u_2(t)) \in (C([0, \infty); L^2))^2$ which satisfies the estimates

$$\left\| u_1(t) + e^{\frac{it}{2m_1}\Delta} \tilde{S}_1(t) u_{1+} \right\|_{L^2} + \left\| u_2(t) + e^{\frac{it}{2m_2}\Delta} \tilde{S}_2(t) u_{2+} \right\|_{L^2} \leq Ct^{-b}$$

for all $t \geq 0$, where $\frac{1}{2} < b < 1$,

$$\tilde{S}_1(t) = \mathcal{F}^{-1} D(m_1) e^{\frac{i}{2}(\theta_0(\xi) + \sqrt{2}|\hat{u}_{1+}(\xi)| \log t)} \mathcal{F},$$

$$\tilde{S}_2(t) = \mathcal{F}^{-1} D(m_2) e^{(\pi - \theta_0(\xi) + 2|\hat{u}_{2+}(\xi)| \log t)} \mathcal{F}$$

and $\theta_0(\xi) = -\arg \hat{u}_{1+}(\xi) + \frac{1}{2} \arg \hat{u}_{2+}(\xi)$.

Remark. Hayashi-Li-Naumkin [1] constructed the modified wave operator with the special assumptions on both amplitude and argument of two scattering states. Theorem 1 shows that we are able to construct the modified wave operator without no restrictions on the argument of two scattering states. We introduce a new angular modification to show the result.

4 Free evolution operator

We introduce decomposition of the free evolution operator.

$$e^{\frac{it}{2m}\Delta} = M_m(t) D_m(t) \mathcal{F} M_m(t),$$

where $M_m(t)$ is the multiplication operator and $D_m(t)$ is the dilation operator defined as follows:

$$M_m(t)\phi(x) = e^{\frac{im|x|^2}{2t}}\phi(x), \quad D_m(t)\phi(x) = \frac{m}{it}\phi\left(\frac{mx}{t}\right).$$

We have formal asymptotics of free Schrödinger equation:

$$e^{\frac{it}{2m}\Delta}\phi \sim M_m(t) D_m(t) \mathcal{F}\phi.$$

5 Outline of the proof

We prove

$$\Phi_1(u) = -i \int_t^\infty e^{\frac{i(t-\tau)}{2m_1}\Delta} (\bar{u}_1 u_2 - \bar{w}_1 w_2) d\tau + \text{h.o.t.}$$

and

$$\Phi_2(u) = -i \int_t^\infty e^{\frac{i(t-\tau)}{2m_2}\Delta} (u_1^2 - (w_1)^2) d\tau + \text{h.o.t.}$$

are contraction mapping on

$$X := \left\{ \phi = (\phi_1, \phi_2) \in (C([T, \infty); L^2))^2 ; \|\phi - w\|_X < \infty \right\}$$

with

$$(w_1, w_2) = (M_{m_1}(t) D_{m_1}(t) \mathcal{F} \tilde{S}_1(t) u_{1+}, M_{m_2}(t) D_{m_2}(t) \mathcal{F} \tilde{S}_2(t) u_{2+})$$

and

$$\|\phi\|_X := \sum_{j=1}^2 \sup_{t \in [T, \infty)} \left(t^{\frac{\beta}{2} + \mu} \|\phi_j(t)\|_{L^2} + t^\mu \| |J_{m_j}|^\beta \phi_j(t) \|_{L^2} \right),$$

where $1 < \beta < \alpha < 2$, $\alpha - \beta > \mu > 0$ and

$$|J_{m_j}|^\beta(t)\phi = e^{\frac{it}{2m_j}\Delta} |x|^\beta e^{-\frac{it}{2m_j}\Delta} \phi.$$

References

- [1] N. Hayashi, C. Li, P. I. Naumkin, *Modified wave operator for a system of nonlinear Schrödinger equations in 2d*, preprint.

Leaf-wise intersections in coisotropic submanifolds

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Introduction

J. Moser defined the leaf-wise intersection which is motivated by perturbation theory of Hamiltonian systems. Perturbation theory is used to find an approximate solution to a problem which cannot be solved exactly.

In this poster, we consider the existence of leaf-wise intersections.

Definitions

(P^{2n}, ω) is called a *symplectic manifold* if ω is a nondegenerate closed 2-form on a $2n$ -dimensional manifold P . A submanifold M of P is said to be *coisotropic* if $(T_p M)^\omega \subset T_p M$ holds for each $p \in M$, and is said to be *Lagrangian* if $(T_p M)^\omega = T_p M$. Here, $(T_p M)^\omega := \{v \in T_p P \mid \omega(v, w) = 0 \ (\forall w \in T_p M)\}$. For a function $H : P \rightarrow \mathbb{R}$, we define the *Hamiltonian vector field* X_H by $i(X_H)\omega = dH$ and denote the flow of X_H by ϕ_H^t . A diffeomorphism $\psi : P \rightarrow P$ is called *Hamiltonian* if $\psi = \phi_H^1$ for some H and we denote the set of Hamiltonian diffeomorphisms by $\text{Ham}(P, \omega)$.

Let M be a coisotropic submanifold of P whose codimension is r . Then $(TM)^\omega \subset TM$ is a completely integrable distribution on M of rank r . Therefore $(TM)^\omega$ defines a foliation on M which is called the *characteristic foliation*. We denote the *leaf* of $(TM)^\omega$ through p by L_p , namely L_p is the maximal r -dimensional submanifold of M through p which is tangent to $(TM)^\omega$. Then $p \in M$ is called a *leaf-wise intersection* of $\psi \in \text{Ham}(P, \omega)$ if $\psi(p) \in L_p$.

Examples

1. $P = \mathbb{R}^{2n}$ with a 2-form $\omega_0 = \sum dx_j \wedge dy_j$ is a symplectic manifold. In this case, the orbits of a Hamiltonian vector field X_H are nothing but the solutions of the Hamilton-Jacobi equations: $\dot{y} = -\frac{\partial H}{\partial x}$, $\dot{x} = \frac{\partial H}{\partial y}$.

2. Let (P, ω, H) be a Hamiltonian system and (P, ω, \tilde{H}) be a perturbed system, namely H and \tilde{H} are very close to each other with their derivatives. We assume that the regular level set $M := H^{-1}(0)$ is compact and the orbits of X_H on M are all periodic with period $T > 0$. If $p \in M$ is a leaf-wise intersection of certain Hamiltonian diffeomorphism, then the orbit of $X_{\tilde{H}}$ through p is periodic. Therefore, the existence of a leaf-wise intersection guarantees the existence of a periodic orbit for the perturbed system.

The first existence theorem

Moser showed the following theorem for the existence of the leaf-wise intersections:

Theorem 1 (Moser, 1978). *Let $(P, \omega = d\lambda)$ be a simply connected exact symplectic manifold and M be a compact coisotropic submanifold of P . If $\psi \in \text{Ham}(P, \omega)$ is C^1 -close to the identity $\text{id}_P : P \rightarrow P$, then ψ has at least $\text{cat}(M)$ leaf-wise intersections. Here, $\text{cat}(M)$ is the Lusternik-Schnirelmann category of M .*

We remark that the simply connectivity assumption can be replaced by the condition that $H^1(P, \mathbb{R}) = 0$.

Main Theorem

We considered the another approach to the existence using a theorem by A. Weinstein and obtained the following:

Theorem 2 (Main theorem). *Let (P, ω) be a symplectic manifold and M be a coisotropic submanifold of P . If $\psi \in \text{Ham}(P, \omega)$ is C^1 -close to the identity $\text{id}_P : P \rightarrow P$, then there exist a closed 1-form Γ on M and an embedding $G : M \rightarrow P \times P$ so that $\text{pr}_1 \circ G(p)$ is a leaf-wise intersection of ψ for each $p \in \text{Zero}(\Gamma)$. In particular, if M is compact and $H^1(M, \mathbb{R}) = 0$, then ψ has at least $\text{cat}(M)$ leaf-wise intersections.*

Actually, it is sufficient that ψ is defined on a neighborhood of M . So the assumptions on the ambient manifold P are not essential. In this meaning, $H^1(M, \mathbb{R}) = 0$ is more rational assumption than $H^1(P, \mathbb{R}) = 0$.

Further results

H. Hofer introduced the norm on the space of compactly supported Hamiltonian diffeomorphisms. For a compactly supported smooth function $H : P \rightarrow \mathbb{R}$, we define $\|H\| := \max H - \min H$. Then, define the norm of a compactly supported Hamiltonian diffeomorphism $\psi \in \text{Ham}_c(P, \omega)$ by $\|\psi\| := \inf \{\|H\| \mid \psi = \phi_H^1\}$.

Theorem 3 (Hofer, 1990). *Let M be a restricted contact type hypersurface in \mathbb{R}^{2n} with the standard symplectic structure. Assume that M bounds an open bounded domain U . Then there exist a constant $c(M) > 0$ such that for any $\psi \in \text{Ham}_c(P, \omega)$ with $\|\psi\| \leq c(M)$ has a leaf-wise intersection.*

Here, a coisotropic submanifold M^{2n-r} is of (restricted) contact type if there exist 1-forms $\alpha_1, \dots, \alpha_r$ on M (on P) such that $d\alpha_j = \omega$ on M (on P) and $\alpha_1 \wedge \dots \wedge \alpha_r \wedge \omega^{n-r} \neq 0$ on M . We note that all hypersurfaces are coisotropic. We remark that $c(M)$ in the theorem is a symplectic capacity.

The other results are listed in the following table.

authors	assumptions on P and M	conclusion
D.L. Dragnev (2008)	$P = \mathbb{R}^{2n}$, M is a compact and contact type coisotropic submanifold	If $\ \psi\ \leq c(M)$, then ψ has a leaf-wise intersection
V.L. Ginzburg (2007)	P is a subcritical Stein manifold, M is a closed connected restricted contact type hypersurface	If $\ \psi\ \leq c(M)$, then ψ has a leaf-wise intersection
P. Albers, U. Frauenfelder (2010)	P is convex at infinity, M is a closed restricted contact type hypersurface	If $\ \psi\ \leq c(M)$, then ψ has $\sum b_i(M, \mathbb{Z}_2)$ leaf-wise intersections
B.Z. Gürel (2010)	P is geometrically bounded and wide, M is restricted contact type coisotropic submanifold	If $\ \psi\ \leq c(M)$, then ψ has a leaf-wise intersection
F. Ziltener (2010)	P is geometrically bounded, M is a closed coisotropic submanifold, the characteristic foliation is a fibration	If $\ \psi\ \leq c(M)$, then ψ has $\sum b_i(M, \mathbb{Z}_2)$ leaf-wise intersections

The origin of low-luminosity AGN/AGN-like activity in red early-type galaxies

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1. Abstract

Primary Aim: We focus on the ionizing mechanism of low-luminosity AGNs(LL-AGNs, especially a class of LINER) in order to construct LL-AGN picture.

Sample: We reconstructed morphologically selected early-type galaxy sample from the SpecObj catalog of the Sloan Digital Sky Survey(SDSS).

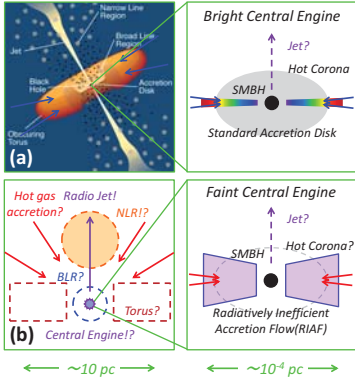
Analyses: We examined emission-line properties and their relations to host stellar mass and $u'-r'$ color.

Results: We obtained two main results: (1)The emission-line ratios of LINERs show the concentrated distribution on the diagnostic diagrams.

(2)The host galaxies of these LINERs always have red color and not always have large stellar mass.

Discussion: The host color suggests that photoionization from old stellar population(possibly post-AGB stars) is promising ionizing mechanism of LINERs and photoionization from LL-AGN is also possible if host stellar age correlates with mass accretion rate onto its super massive black hole(SMBH).

2. Introduction



Active Galactic Nuclei(AGN) pictures.

Bright AGN(Fig.a): When matter accrete onto SMBH, they release their gravitational energy. This energy is transformed efficiently into radiative energy in accretion disk. This radiation ionizes ambient gases. As a result, ionized gases radiate prominent emission lines.

Faint AGN(Fig.b): When accretion rate decreases, accretion disk state changes into RIAF. Therefore, central region and emission lines become faint. However, faint AGN picture is still under discussion because these weak activities can also be explained by many other models.

Fig.a: Urry & Padovani 1995, PASP 107, 803

3. Sample

Flux limited sample of nearby bright, morphologically selected, 34497 ETGs.

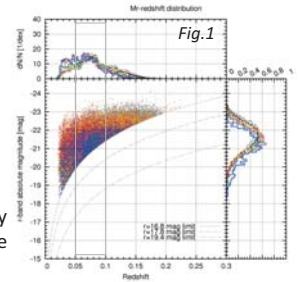
Data set:

- SDSS DR7 SpecObj catalog
- the MPA-JHU DR7 release of spectrum measurements

Selection criteria:

fracDeV_g', r' , i' > 0.95,
Spectroscopic S/N_g', r' , i' > 10,
 $0.05 < z < 0.1$, $r' < 16.8$.

Fig.1: The distribution of our galaxy sample in r -band absolute magnitude vs. redshift plane.



4. Analyses & Results

4-1. Emission-line diagnostic diagrams

$H\beta$, $[OIII]\lambda 5007$, $H\alpha$, $[NII]\lambda 6584 > 3\sigma \rightarrow$ Active ETGs,
 $1 \sim 3$ these emission lines $> 3\sigma \rightarrow$ Semi-Active ETGs,
 $H\beta$, $[OIII]\lambda 5007$, $H\alpha$, $[NII]\lambda 6584 < 3\sigma \rightarrow$ Quiescent ETGs.

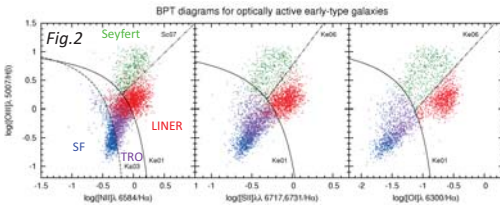


Fig.2: Result of our classification. From fig.2, we recognized clump of the distribution in LINER region. To investigate correlation to other properties of these LINERs, we also plotted the diagrams with the coloring according to its host stellar mass or color.

Ka03: Kauffmann et al. 2003, MNRAS 346, 1055, Ke01: Kewley et al. 2001, ApJ 556, 121,
Ke06: Kewley et al. 2006, MNRAS 372, 961, Sc07: Schawinski et al. 2007, MNRAS 382, 1415
SF: Star Forming
TRO: Transition Region Object
LINER: Low-Ionization Emission-Line Region

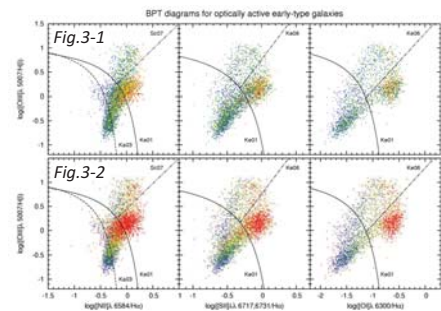


Fig.3: The diagnostic diagrams with the coloring defined above. These diagrams show that massive or red active ETGs tend to concentrate into the LINER region. Because host stellar mass correlates with its color, however, it is possible that we see just the same ETGs. So next, we plotted the active ETGs separately.

After classification:
830 SF
1576 TRO
637 Seyfert
1897 LINER
12444 Semi-Active
17113 Quiescent

Color definition

$\log(M_*/M_\odot) < 10.4$,
 $10.4 < \log(M_*/M_\odot) < 10.7$,
 $10.7 < \log(M_*/M_\odot) < 11.0$,
 $11.0 < \log(M_*/M_\odot) < 11.3$,
 $11.3 < \log(M_*/M_\odot)$
 $u'-r' < 1.9$,
 $1.9 < u'-r' < 2.2$,
 $2.2 < u'-r' < 2.5$,
 $2.5 < u'-r' < 2.8$,
 $2.8 < u'-r'$

4-2. Color-Mass & Diagnostic diagrams

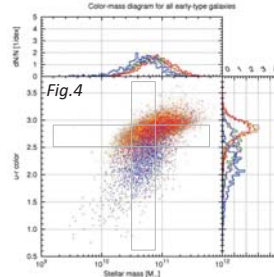
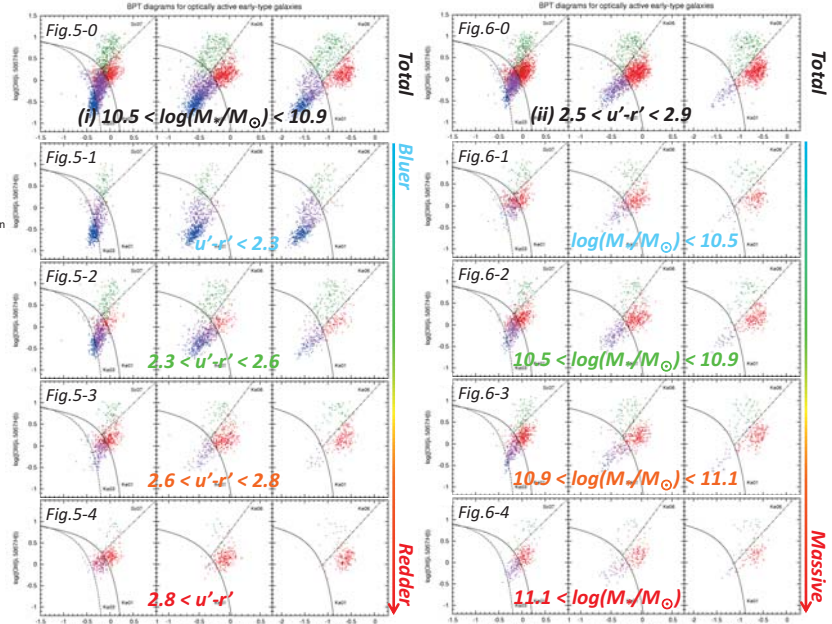


Fig.4: Color-mass diagram for our classified ETGs. We intended to show our separation range in these panels(gray square).

Case (i) : $10.5 < \log(M_*/M_\odot) < 10.9$

Case (ii) : $2.5 < u'-r' < 2.9$

The diagnostic diagrams for each subsample are showed in figs.5 – 6. [Note that figs.5-0 & 6-0 show total plots of the case (i) & (ii)]



(i) In figs.5-1 – 5-4, each subsample shows very different distributions in the sense that redder subsample concentrates into the LINER region.

(ii) In figs.6-1 – 6-4, this distribution dose not vary across subsamples.

5. Discussion

A number of the ionizing mechanisms reproducing LINER(-like) emission-line properties have been proposed. These mechanisms fall roughly into four categories: (A)photoionization by LL-AGN, (B)fast shock, (C)photoionization by post-AGB stars(pAGBs), (D)interaction between warm and hot inter-stellar medium. If we adopt (A), we speculate that (1)as the old stellar population dominates in host galaxy, mass accretion rate onto its SMBH decreases, (2)accretion state varies from the standard accretion disk to radiatively inefficient accretion flow(RIAF), (3)the ionization parameter decreases with maintaining the hardness of ionizing spectrum, (4)as a result, such galaxies become LINERs. If we adopt (C), it suggests that the pAGBs spreading over host galaxy provide ionizing photons. In (B) and (D), because these can not simply explain the host properties presented here, we think it is not likely that these mechanisms are primary in the clump LINERs.

Path Integral Analysis of Bianchi type I spacetime in Loop Quantum Cosmology

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Loop quantum gravity (LQG) is the best candidate for quantum gravity. Cosmology based on LQG is called Loop quantum cosmology (LQC).

Although the full analysis is very complex even in toy model. we can use the path integral formulation to investigate the brief behavior.

In this presentation, path integral analysis of Bianchi type I spacetime with a massless scalar field and cosmological constant will be reported.

Bianchi type I spacetime

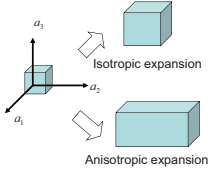
Bianchi type I spacetime is homogeneous but anisotropic model. There are three directional scale factors. However their values have no physical meaning directly, only ratios $a_i(\tau)/a_i(\tau_0)$ have.

$$ds^2 = a_1^2 dx_1^2 + a_2^2 dx_2^2 + a_3^2 dx_3^2$$

$$a_i = ae^{\theta_i}, a = (a_1 a_2 a_3)^{1/3} \quad H_i = \frac{\dot{a}_i}{a_i} \quad \text{Directional Hubble parameter (directional expansion rate)}$$

Note that $\dot{\theta}_i \propto \frac{1}{a}$ In classical general relativity

Anisotropy of spacetime is defined as $\Sigma = \frac{1}{18}[(H_1 - H_2)^2 + (H_2 - H_3)^2 + (H_3 - H_1)^2]$
 $= \frac{1}{6}(\dot{\theta}_1^2 + \dot{\theta}_2^2 + \dot{\theta}_3^2) \propto \frac{1}{a^6}$ In classical general relativity



Friedmann equation is written as

$$H^2 = \frac{8\pi G}{3} \rho + \Sigma + \frac{\Lambda}{3}$$

Effective action of Bianchi I model

We investigate Bianchi I spacetime with a massless scalar field ϕ and cosmological constant Λ . The effective action is given as;

$$S_{\text{Bianchi I}} = \int d\tau \left(p_\phi \dot{\phi} - \frac{\beta_1 \dot{\lambda}_1}{2} - \frac{\beta_2 \dot{\lambda}_2}{2} - \frac{b}{2} \dot{v} + \alpha \left[p_\phi^2 + \frac{\pi G}{4l^2} [v^2 P(\lambda, v, \beta, b) + 2\hbar v Q(\lambda, v, \beta)] + \pi G \Lambda v^2 \right] \right)$$

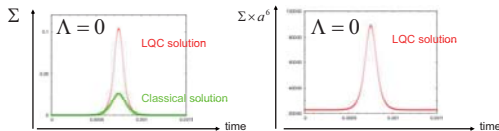
$\lambda_i \propto \sqrt{a_i a_k}$ (i,j,k cyclic) p_ϕ, β_i and b are momenta to ϕ, λ_i and volume respectively

$l = \gamma \sqrt{4\pi G \hbar^2}$, γ is Barbero-Immirzi parameter ~ 0.2735

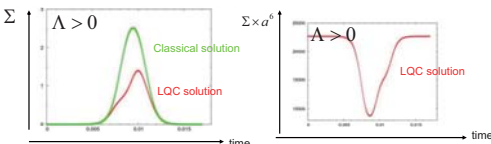
$$P, Q = \cos \left[\frac{l}{\sqrt{v(v+2\hbar)}} (\beta_1 \lambda_1 - \beta_2 \lambda_2) \right] + \dots$$

We can obtain first order approximation equations from this action. These equations reproduce the general relativity in large scale, but there is non-trivial departure from classical theory in small scale.

We can show that $\beta_i \lambda_i$ are constant and they represent anisotropy. Therefore if we set $\beta_i = 0$, this model becomes the isotropic one.



The anisotropy Σ is larger than classical one in small scale. However in $\Lambda > 0$ case have contrary result. The anisotropy in LQC is smaller than classical one.



This is because the anisotropy Σ has upper limit $\Sigma_{\text{max}} \sim 1.85$ in LQC. Therefore we summarize as

- LQC anisotropy is larger than classical one if classical anisotropy is smaller than $\Sigma_{\text{max}} \sim 1.85$
- LQC anisotropy is smaller than classical one if classical anisotropy is larger than $\Sigma_{\text{max}} \sim 1.85$
- Note that in both case the anisotropy in large scale is conserved under the big bounce.

Loop Quantum Cosmology

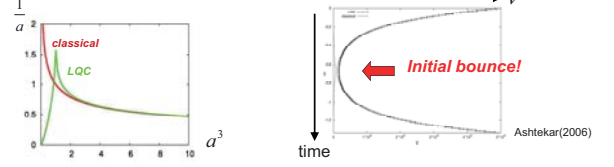
Cosmology based on Loop Quantum Gravity (LQG) is called Loop Quantum Cosmology (LQC).

- LQC has NOT *differential* equation but *difference* equation. This is caused by discreteness of geometrical values.

$$(\hat{C}_{\text{grav}} + \hat{C}_{\text{mat}}) \Psi(v) = f_+(v) \Psi(v+4) + f_0(v) \Psi(v) + f_-(v-4) \Psi(v-4) + \hat{C}_{\text{mat}} \Psi(v) = 0$$

v : volume of the Universe

- Expectation value of volume has initial bounce. This means that the quantum effect causes the repulsive force to avoid the singularity.



Path integral formulation

Although we play with simple toy models, whether analytic or numerical, full Hamiltonian analysis of LQC is very complicated.

Ashtekar followed the Feynman's original path integral formulation. We can obtain *effective action* which contain *quantum correction*.

$$\langle v_f, \phi_f | v_i, \phi_i \rangle = \delta_{v_f, v_i} \delta(\phi_f - \phi_i) \quad A(v_f, \phi_f; v_i, \phi_i) \equiv \int d\alpha \langle v_f, \phi_f | e^{i\alpha \hat{C}} | v_i, \phi_i \rangle$$

$$(\Phi_{\text{phys}}, \Psi_{\text{phys}}) \equiv \langle \Phi_{\text{kin}} | \int d\alpha e^{i\alpha \hat{C}} | \Psi_{\text{kin}} \rangle = \sum_{v, \phi} \int d\phi d\phi' \overline{\Phi}_{\text{kin}}(v, \phi) A(v, \phi; v', \phi') \Psi_{\text{kin}}(v', \phi')$$

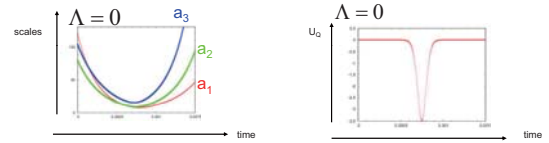
$$\langle v_f, \phi_f | e^{i\alpha \hat{C}} | v_i, \phi_i \rangle = \sum_{v_N, \phi_N} \int d\phi_N d\phi_{N-1} \dots d\phi_1 \langle v_f, \phi_f | e^{i\alpha \hat{C}} | v_N, \phi_N \rangle \dots \langle v_1, \phi_1 | e^{i\alpha \hat{C}} | v_i, \phi_i \rangle$$

$$A(v_f, \phi_f; v_i, \phi_i) = \int d\alpha [Dv(\tau)] [D\phi(\tau)] [Dp(\tau)] e^{iS}$$

$$S = \int d\tau \left(p \dot{\phi} - \frac{1}{2} b \dot{v} - \alpha (p^2 - 3\pi G v^2 \sin^2 b) \right) \quad \text{Flat FRW + massless scalar field model (Ashtekar's original paper)}$$

Results

We compute the solution of first order approximation equations. At first we consider $\Lambda = 0$ case.



Scale factors are bouncing in small scale as isotropic one. We can write the quantum effect $U_0 = H^2 - \frac{3\pi G}{4} \rho - \Sigma - \frac{\Lambda}{3}$ i.e. departure from classical Friedmann eq. In large scale $U_0 \sim 0$ but becomes large with negative sign in small scale. This shows that quantum effect causes the repulsive force to avoid the initial singularity.

$\dot{\theta}_i \propto \frac{1}{a^3}$ In classical theory. However this relation is invalid because of quantum effect. Also

$$\Sigma = \frac{1}{6}(\dot{\theta}_1^2 + \dot{\theta}_2^2 + \dot{\theta}_3^2) \propto \frac{1}{a^6} \quad \text{Is invalid in LQC.}$$

Summary

- The effective action of Bianchi I spacetime with a massless scalar field and cosmological constant is constructed. We obtain approximation equations which contain quantum correction. And these equation has isotropic solution as special case.
- In large scale, equations reproduce the classical general relativity. Even if there is anisotropy, the quantum effect cause the repulsive force in small scale to avoid the Big Bang singularity as we expected.
- The anisotropy has interesting behavior in small scale. Because the anisotropy has upper limit in LQC, If classical anisotropy is larger (smaller) than this upper limit, LQC anisotropy is smaller (larger) than classical one. The anisotropy in large scale is conserved under the big bounce. Note that results are depending on matter configuration. If we consider matter content which has anisotropic stress (e.g. magnetic field), the results maybe change.

Consistency Relation for Multi Field Inflation Scenario

Naonori Sugiyama

Abstract

While it is well-known that the universe began from a very hot and dense state called "Big Bang", many cosmologists now think that there was a period even before the "Big Bang". It is called "inflation". In this "inflation period", the universe has expanded more rapidly than it does today, and this is the reason for its name "inflation". Unfortunately, we have not yet obtained a satisfactory physical description of inflation because we do not have cosmological observation accurate enough to reach it. Indeed, lack of such observations allows theorists invent many inflation models different from each other when comparing those in various aspects. In short, "The number of inflation models is being inflated".

In this work, we derived a quite general inequality among observable quantities by considering a concept called "Non-Gaussianity". This inequality should be satisfied by almost all inflation models considered today. If future observation suggests that our inequality does not hold in our universe, our inequality can rule out almost all models subject to it.

The universe began from "Big Bang"!!

Big Bang



A huge explosion



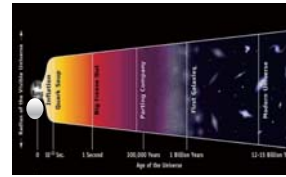
A much hot and dense state

"Big Bang" is not a huge explosion. It is a much hotter and denser state of the universe than today. In the early universe, there are just elemental matters, for example electron, neutron, and photon. The universe have been expanding and cooling ever now.

The universe began from "Big Bang"!!

The universe was a hotter and denser state called "Big Bang"!!

There was a period called "inflation" before "Big Bang"!!



Many cosmologists think that the universe was expanding much more rapidly before "Big Bang" than it does today.

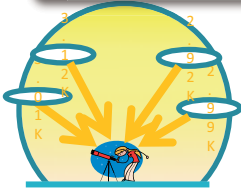
We call the period "inflation".

Inflation theory is just a theoretical demand, it is not yet verified from observation. Furthermore nobody understand the detail of physics during inflation period.

The verification of inflation theory is one of the most important topic in cosmology.

How do we verify the inflation theory?

Cosmic Microwave Background radiation (CMB)



There is the light called "Cosmic Microwave Background radiation" in the universe. Its temperature is 2.725 K (Kelvin). It is the oldest light that we can observe. To be precise, the temperature of the CMB is different in each direction it comes. The distribution for the difference of the temperature is nearly "Gaussian distribution". Recently, it has been found that CMB has a little departure from Gaussian distribution. It is called "Non-Gaussianity of the perturbation of the CMB temperature". Many cosmologists think that inflation scenario explain how the Non-Gaussianity is created. So we have information of the inflation period by researching the CMB Non-Gaussianity.

In order to characterize amount of the Non-Gaussianity, we often use the below parameters, called "Non-linear parameter",

$$\tau_{NL} \quad f_{NL}$$

The way of theoretical research for inflation

We have not yet understand inflation physics in detail. Each researcher makes their own models for inflation and predict s different non-linear parameters.

There are inflation models which predict a small non-linear parameter: $f_{NL} \approx O(10^{-2})$

While, a large non-linear parameter is also predicted in other models: $f_{NL} \approx O(10)$

The observational value of non-Gaussian parameters

$$f_{NL} = 32 \pm 21(65\%CL)$$

τ_{NL} is not yet observed.

We do not have observational values enough to constrain theoretical Inflation models. So we can not rule out almost all inflation models. As the result, the number of Inflation models have been increasing.

It is difficult to verify huge number of inflation models one by one. We should have a more general discussion.

$$\tau_{NL} \geq \left(\frac{6}{5} f_{NL}\right)^2$$

We have discovered above inequality between non-linear parameters.

This inequality is satisfied in a large class of inflation models called "multi-field inflation".

If the new experimental data such as Planck (which is a satellite for observation of CMB) detect f_{NL} but do not see τ_{NL} large enough to satisfy the above inequality, then almost all inflation models may be ruled out.

Conclusion

We have discovered the general inequality (A) between non-linear parameters. The most interesting case would be the observation of a complete violation of the inequality, i.e.,

$$\tau_{NL} < \left(\frac{6}{5} f_{NL}\right)^2$$

which implies that inflation cannot be responsible for generating the observed Non-Gaussianity. We may not be so far away from testing this prediction. Our result provide a strong science case for measuring the local-form tri-spectrum of the cosmic microwave background.

The star formation in the SSA22 protocluster at $z=3.09$

2012/2/20-22 GCOE symposium

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1. Overview; K-selected galaxies in the SSA22 protocluster

The SSA22 $z=3.09$ protocluster is a known high density region at high redshift which identified with the density excess of the Ly α Emitters(LAEs), Lyman break Galaxies(LBGs) at $z\sim 3.1$.

In this work, we selected the candidate of the protocluster galaxies based on stellar mass with MOIRCS JHK band deep imagings at the highest density region of LAEs in the SSA22 protocluster. And then analysed the star formation and stellar populations in the protocluster using Spitzer IRAC, MIPS 24 μ m (Webb+2008). We selected the galaxies with $K_{AB} < 2.4$ at $z_{phot} = 2.6-3.6$ by using photometric redshift estimated from Spectral Energy Distribution(SED) fitting with UBVRI'zJHK, IRAC 3.6,4.5, 5.8, 8.0 μ m. The photometric redshift error at $z\sim 3$ is ~ 0.5 . We selected the galaxies with stellar mass $\sim 10^{9-11} M_{sun}$.

The sky distribution of the K-selected galaxies is Fig.1. They are concentrated around the density peak of the LAEs at $z=3.09$. The surface number density of the K-selected galaxies is 1.6 times larger than that in the GOODS-North field at same redshift range(Fig. 2).

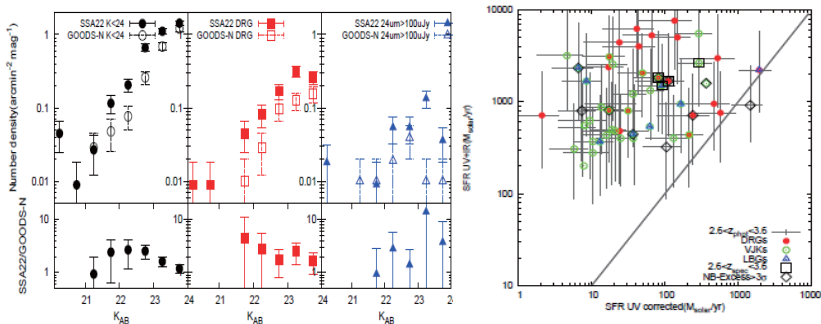


Figure. 2; The comparison of the differential surface number density of the K-selected galaxies(left), those classified as DRGs(middle), MIPS24 μ m detected (right) in the SSA22 protocluster(filled) and in the GOODS-North field(Kajisawa et al.2009) (unfilled). Bottom panels are the ratio of the number density in the SSA22 to that in the GOODS-N.

Figure.3; The SFR_{UV-IR} v.s. $SFR_{UVcorrected}$ of the 24 μ m sources. K-selected galaxies selected as DRGs(red), VJKs(green, BzK galaxies at $z\sim 3$), LBGs(blue) are marked. Black squares are those at $z_{spec}=2.6-3.6$, black diamonds are those with NB-excess $>3\sigma$.

3. The passive galaxies

Furthermore, we found the passively evolving galaxies in the protocluster.

Though, the detection limit of 24 μ m is too large, we used rest-UV to NIR color to see the properties of the galaxies. Fig.4 is the I'-K v.s. K-4.5 μ m color diagram. 24 μ m detected galaxies have color consisted with dusty starburst. On the other hand, there are the galaxies which have color like the single burst >0.5 Gyr, those would be passively evolving galaxies.

Top panel of the Fig. 4 is the K-4.5 μ m color distribution of the DRGs in the SSA22 protocluster and in the GOODS-North field at $z=2.6-3.6$. There are the excess of the K-4.5 μ m bluer DRGs, or passive galaxies in the protocluster. Sky distribution of the passive galaxies(in Fig.1) show strong clustering around the highest density region of the LAEs. The x-ray detected, 24 μ m detected galaxies show similar sky distributions.

4.Discussion and Future works

There are the density excess of the K-selected galaxies detected with 24 μ m. They have very high SFR and they would rapidly evolve into the massive ellipticals in the center of the present cluster. Furthermore the sky and color distribution of the X-ray detected galaxies suggest the correlation of the nuclear activities and the formation of the massive galaxies.

Are the passive galaxies already like the local massive ellipticals? The stellar mass of these are as large as $10^{10.5-11} M_{sun}$. Fig.5 are the resolved image of 24 μ m source(top) and the passive galaxy selected with I'-K >3.0 and K-4.5 μ m >0.5 (bottom). At least this passive galaxy is massive spheroid like system.

Thus, at the center of the $z=3.1$ protocluster, the massive galaxies are dramatically being formed but some mature systems may already exist.

Future works, the NIR spectroscopies of the K-selected galaxies are required to confirm their redshift, though there are large uncertainties about photometric redshifts.

Morphological analysis of these are also important to understand the formation and evolution of the massive ellipticals.

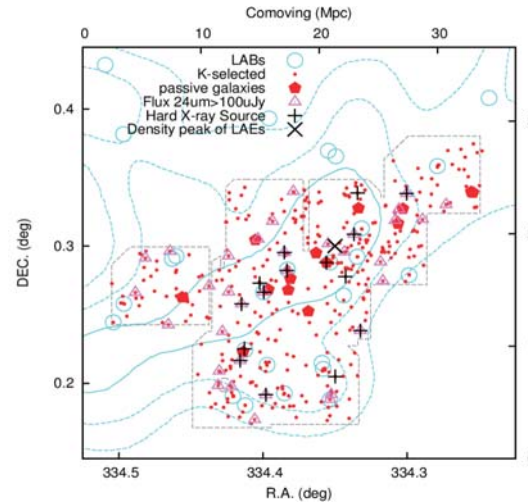


Figure. 1; The sky distribution of the K-selected galaxies ($K_{AB} < 2.4$ and $2.6 < z_{phot} < 3.6$) (redcircles). The large X is the density peak of the LAEs. The light blue circles are Ly α Blobs(LABs). Big red pentagons are the passive galaxies selected with I'-K >3.0 and K-4.5 μ m <0.5 . The 24 μ m detected(pink triangles), X-ray detected(Lehmer+2009, black crosses) are marked.

2. The excess obscured starbursts in the protocluster

There are further density excess of some K-selected galaxies; which selected as Distant Red Galaxies(DRGs; $J-K_{AB} > 1.4$) are 2.2 times, which detected with 24 μ m are 3.4 times numerous (Fig.2).

Fig.3 is the SFR_{UV-IR} and $SFR_{UVcorrected}$ of the 24 μ m detected galaxies. SFR_{UV-IR} are estimated from unobscured UV light and dust obscured star formation re-emitted at IR. $SFR_{UVcorrected}$ are estimated from the UV light corrected of the extinction.

They have $SFR \sim 1000 M_{sun}/yr$. Most of them are too faint at rest-UV to be selected as LBGs, LAEs, but be selected as DRGs.

Thus NIR survey have revealed not only the stellar mass but also the obscured starburst in the protocluster.

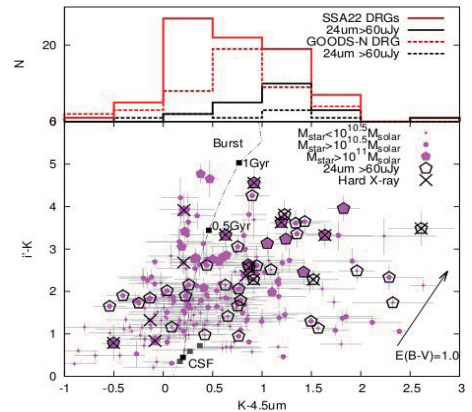


Figure.4; The bottom panel is I'-K v.s. K-4.5 μ m color distribution of the K-selected galaxies(pink pentagons). 24 μ m detected galaxies are marked with black pentagons. Black dashed line is the age evolution track of the single burst star formation model at $z\sim 3$. The grey line is same but of the constant star formation model.

Top panel of Figure.5 is the K-4.5 μ m color distribution of the DRGs(red) at $2.6 < z_{phot} < 3.6$ and those detected with 24 μ m (black) in the SSA22 proto-cluster(solid lines) and in the GOODS-North field (dashed lines).

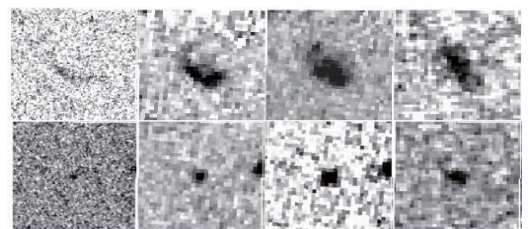


Fig.5 The HST ACS F814(PI Chapman ID 10405), WFC3 F110W, F160W(PI Siana ID 11636), and MOIRCS K images of the K-selected galaxies.

Top; 24 μ m detected galaxy. Bottom; passive galaxy with I'-K >3 and K-4.5 μ m <0.5

Quasinormal modes of charged anti-de Sitter black holes

Nami Uchikata and Shijun Yoshida (*Astronomical Institute, Tohoku University*)

Abstract

It is well known that rotating and/or charged black holes in asymptotically flat spacetime become unstable via superradiant scattering. Superradiance is a phenomenon that an impinging wave can be amplified by a rotating and/or charged black hole when the wave satisfies a certain condition.

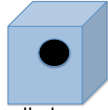
For the black holes in asymptotically anti-de Sitter (AdS) spacetime, it has been known that the rotating AdS black holes also become unstable if their quasi-normal modes satisfy the superradiance condition when the black holes' radii are sufficiently smaller than the AdS scale. However, it is not well studied for the charged AdS black holes. Thus, we have investigated the quasinormal modes of the charged AdS black holes analytically and numerically in the small black hole limit, and showed that the charged AdS black holes become unstable against scalar perturbations if their quasinormal modes satisfy the superradiance condition.

Introduction

[What is an anti-de Sitter (AdS) black hole ?]

AdS is the spacetime whose boundary behaves as an effective wall at spatial infinity. We refer to this characteristic size of the AdS spacetime as an AdS scale (L).

AdS black holes are black holes which exist in such spacetime.



[What are quasinormal modes?]

Every objects have their characteristic oscillations, called normal modes. For black holes, the oscillations usually decay (or amplify) with time, so we say *quasinormal* modes.

(Vishveshwara 1970; Chandrasekhar & Detweiler 1975, Leaver 1985)

Since the AdS boundary behaves as a wall, the wave scattered by the black hole will reflect at the AdS boundary. Thus, the waves having decaying (growing) modes are stable (unstable). Instability occurs only for rotating and/or charged black holes. And this unstable phenomenon is called superradiance.

(Zel'dovich 1971; Misner 1972; Starobinsky 1973)

[Example]

The quasinormal modes of rotating AdS black holes have unstable modes when their radii are sufficiently smaller than the AdS scale. (Cardoso & Dias 2004; Cardoso et al. 2004; Cardoso et al. 2006; Uchikata et al. 2009)

We focus on the quasinormal modes of charged AdS black holes and investigate whether the system is stable or not.

Calculation methods

We consider the perturbations of the charged scalar field on the charged AdS black holes.

[Metric of charged AdS black holes]

$$ds^2 = -f(r)dt^2 + f(r)^{-1}dr^2 + r^2(d\theta^2 + \sin^2\theta d\phi^2)$$

$$f(r) = \frac{r^2}{L^2} + 1 - \frac{2M}{r} + \frac{Q^2}{r^2}$$

M : the black hole mass Q : the charge of the black hole

L^2 : the inverse of a cosmological constant (Λ)

[Equations of perturbations]

Quasinormal modes have complex ω . If we decompose the wave function Φ as, $\Phi(r, \theta, \phi, t) = R(r)Y_{\ell m}(\theta, \phi)e^{-i\omega t}$,

then $\Phi \propto \exp(-i \operatorname{Re}(\omega) t) \exp(-\operatorname{Im}(\omega) t)$. ($Y_{\ell m}$ is the spherical harmonics.)

• $\operatorname{Re}(\omega)$ represents frequency of the oscillation.

• $\operatorname{Im}(\omega)$ represents the damping timescale of the oscillation.

Thus, if $\operatorname{Im}(\omega) > 0$, Φ grows exponentially with time.

→ The wave is amplified as it scattered off by the black hole. Thus the system is **unstable**.

• the basic equation (radial part)

$$\frac{d}{dr} \left(r^2 f(r) \frac{dR(r)}{dr} \right) + \left[\frac{(\omega r - eQ)^2}{f(r)} - \ell(\ell+1) \right] R(r) = 0$$

e : the charge of the scalar field

Results

[Analytical calculation]

We find that

$$\operatorname{Im}(\omega) \propto - \left(\operatorname{Re}(\omega) - e \frac{Q}{r_+} \right),$$

for small black holes. (The black hole radius, r_+ , is sufficiently smaller than the AdS scale, $r_+ \ll L$.) Thus, if the condition $\omega - eQ/r_+ < 0$ is satisfied, the black hole becomes unstable. The instability condition can also be written as

$$\frac{Q}{Q_c} > \frac{3}{eL},$$

for the fundamental mode, where Q_c is the maximum of Q .

[Numerical calculation]

It is verified $\operatorname{Im}(\omega) > 0$ when $\omega - eQ/r_+ < 0$ for small black holes.

Fig.1

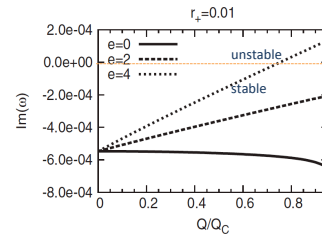
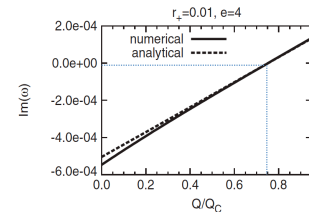


Fig.2



(Uchikata & Yoshida 2011)

Both figures are the results of the imaginary part of the frequency for the black hole with its radius $r_+ = 0.01L$. Fig.1 shows that for the scalar field with $eL < 4$, the black hole is stable but for the scalar field with $eL = 4$, the black hole becomes unstable when the black hole charge, Q , is larger than $0.75Q_c$. Fig.2 shows that numerical results agree with analytical calculations.

Conclusions

- We have studied the scalar perturbations of the charged AdS black holes whose radii are sufficiently smaller than the AdS scale.
- From both analytical and numerical calculations, we have showed that the instability condition can be expressed as $\omega - eQ/r_+ < 0$, or

$$\frac{Q}{Q_c} > \frac{3}{eL}.$$

For example, when the scalar field has a charge $eL = 4$, the black hole becomes unstable when the black hole charge is $Q > 0.75Q_c$.

- It is also interesting to examine whether this instability condition stands for the large AdS black holes, i.e. for black holes having $r_+ > L$. This will be the future problem.

Scanning tunneling microscopy of electronic properties of bulk and layered MoS₂

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¹Tohoku University, ²JST-ERATO



INTRODUCTION

Reaching to dimensional limit of the metal-oxide semiconductors, finding other alternatives is crucial. The suitable substitute should diminish the size of the semiconductor and possess the required characteristics of the semiconductors such as band gap.

Low dimensional electron systems

Low dimensional electron systems can reduce the size of the electronic circuits significantly. The well known two dimensional system, graphene, could be a appropriate choice, however it has no intrinsic band gap which is necessary for logic devices. However a single layer of molybdenum disulfide (MoS₂) has an intrinsic band gap, $E_g=1.8\text{eV}$. This enabled to fabricate MoS₂-based logic device recently[1].

A single crystal of MoS₂ (bulk), likewise graphite, is composed of vertically stack layers, held by weak Van der Waals force.

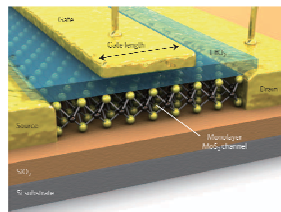
Molybdenum disulfide (Bulk)

- Forms in nature as mineral ore in the form of bulk.
- The appearance is black solid (Similar to graphite).
- Weak van der Waals interaction between the layers.
- Used as lubricant.
- Large indirect band gap, $E_g=1.2\text{eV}$.

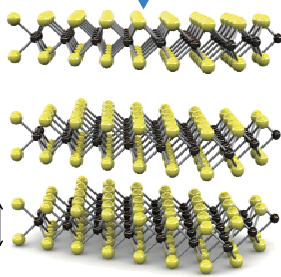
Molybdenum disulfide (single layer)

- Single layer has the thickness of 6.5\AA .
- Direct band gap semiconductor with $E_g=1.8\text{eV}$.

In this work, we observed the surface of the bulk MoS₂ and visualized the distinct step structures using scanning tunneling microscopy (STM). We then examined the band gap of the bulk MoS₂ by scanning tunneling spectroscopy (STS).



Fabrication of MoS₂ monolayer transistors [1].



EXPERIMENT

- Omicron STM in ambient condition.
- Active damping system for mechanical isolation.
- Mechanically sharpened Ir/Pt tip.
- The modulation frequency and the modulation voltage are 1.7KHz and is 15mV for STS measurements.

RESULTS

STM and STS results on bulk MoS₂.

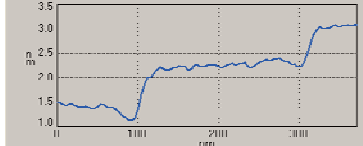
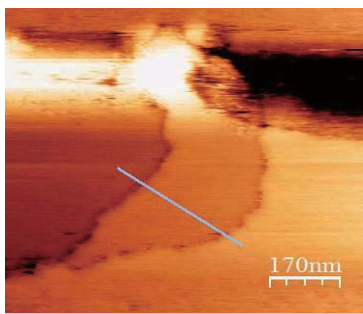


Fig.1: 850x850nm² STM image of bulk MoS₂.

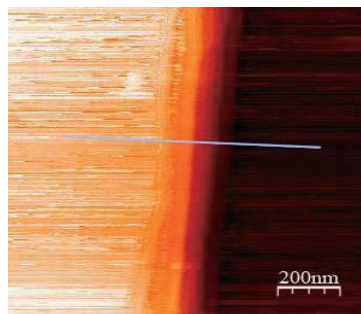


Fig.3: 1x1µm² STM image of bulk MoS₂.

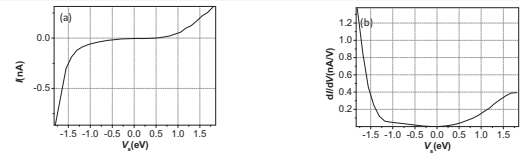


Fig.2: Average I - V and dI/dV curves detected on the area shown in figure 1.

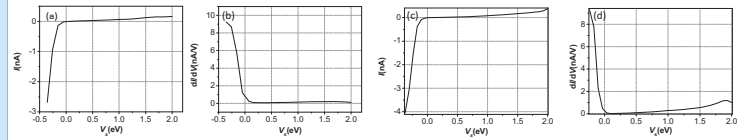


Fig.4: Average I - V and dI/dV curves detected on the area shown in figure 3.

RESULTS

In figure 1 we applied positive sample bias voltage, $V_s=+6.0\text{V}$ and tunnel current, $I_t=0.1\text{nA}$. The bottom panel of figure 1 shows the corrugation across the light blue line from which the step height is estimated to be about 6.5\AA . This is the thickness of a single layer of MoS₂.

Figure 2 (a) shows the average I - V curve taken on the different area shown in figure 1. I_t varies between -0.9nA and 0.3nA while V_s ranges between -1.8V and 1.8V .

In the figure 3, the top panel shows the STM image is taken by applying negative sample bias voltage, $V_s=-0.1\text{V}$ and tunnel current $I_t=0.04\text{nA}$. The bottom panel shows a 7nm thick layer of MoS₂ across the light blue profile line. This thick layer is composed of approximately 10 single layers of MoS₂.

Figure 4 shows the I - V and dI/dV curves detected on the area shown on figure 3 with the set points $V_s=+1.5\text{V}$ and $I_t=0.1\text{nA}$ for (a),(b) and $I_t=0.2\text{nA}$ for (c),(d).

The I - V curves in figure 2 and 4 are asymmetric. It seems in negative sample bias regime the absolute value of I_t increases drastically unlike the positive V_s regime.

DISCUSSION

In figure 2_(b) the dI/dV curve shows a band gap, $E_g=0.7\text{eV}$ which is close to the calculated indirect band gap $E_g=1.06\text{eV}$ [2].

A shift in offset of I - V curves can be caused by tip oxidation or tip contamination. Since the Ir/Pt tip is robust against oxidation it is more likely that an object was attached to the tip apex. It is reported that the a likely tip contaminant is sulfur because of an excess of S at the MoS₂ cleavage surface[4]. However this contaminated apex can lead to clear and correct topographic image but not necessarily a real LDOS.

The curves in figure 4 (b) and (d) in positive V_s regime show different behavior. For $V_s > 0.0\text{eV}$, curve (b) does not change so much, but curve (d) starts to raise for $V_s > 0.5\text{eV}$. This behavior should be canceled out after normalizing by dividing them by I - V curves, however this is not our case (the normalized curves are not shown here). The origin of this behavior is not clear but one possible reason can be tip contamination which can also lead to energy shift as mentioned before.

CONCLUSION

- The topographic image of single layer and multilayer of bulk MoS₂ was obtained in ambient condition.
- A band gap, $E_g=0.7\text{eV}$ is detected for bulk MoS₂.
- A shift in offset of I - V curves and the asymmetric behavior of dI/dV curves can be due to contaminated tip.

PROSPECTIVE

- Detecting a few layers as well as single layer of MoS₂ using atomic force microscope (AFM) (fig.5)

- Studying the electronic properties of bulk and single layer of MoS₂ at low temperature in the presence of high magnetic field using STM/STS.

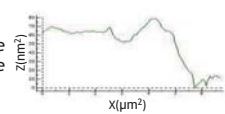


Fig.5: The AFM image of MoS₂ flakes. The left panel shows the corrugated surface across the blue line.

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KamLAND and MoGURA

- ▶ KamLAND is a low energy neutrino and anti-neutrino detector
 - ▷ 1000ton of liquid scintillator in $\phi 13\text{m}$ sphere
 - ▷ ultra clean and low background
 - ▷ converted to a Xenon double beta decay searcher as *KamLAND-Zen*
 - ▷ major challenge for DAQ:
 - energetic cosmic μ overwhelms DAQ with exceeding events
- ▶ MoGURA: Module for General-Use Rapid-Application
 - ▷ new DAQ system operating in parallel for KamLAND
 - ▷ integrates 1GSPS flash ADC and Spartan 3E FPGA from Xilinx
 - ▷ large ring buffer holding 10us waveform

MoGURA System Overview

trigger board



Trigger Logic FPGA: Spartan 3E x 1
Trigger System FPGA: Spartan 3E x 1

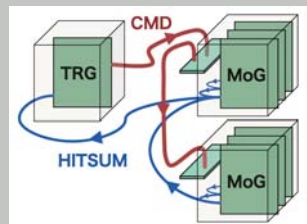
main board



Front End FPGA: Spartan 3E x 6
System FPGA: Spartan 3 x 1
User FPGA: Virtex 2 Pro x 1

- ▶ FPGA: field-programmable gate array
 - ▷ integrated circuit is loaded from flash memory when power on
 - ▷ flash memory are reprogrammable flexibility in developing and debugging
- ▶ Features: to catch the most useful data
 - ▷ trigger system
 - ▷ zero suppression
 - ▷ mode switch by occupancy of buffer normal, almost full, full

illustration of the whole



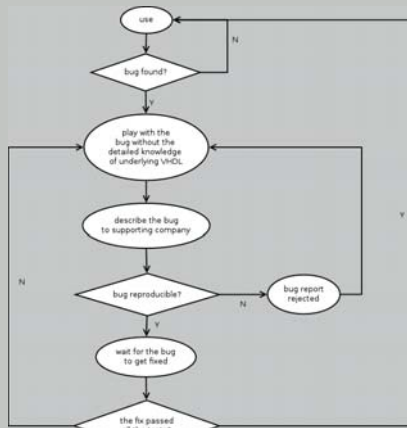
Development and Debug with VHDL

- ▶ VHDL: VHSIC hardware description language
 - ▷ VHSIC: very-high-speed integrated circuits
 - ▷ call for dedicated hardware skills to master (e.g. hardware constraints)
 - ▷ too deep learning curve to be an auxiliary tool for physics
 - ▷ present systems (front end FPGA, system FPGA) are complex

Problems in Present Working Flow

- ▶ it is hard to describe and reproduce the bug
 - ▷ without understanding VHDL
 - ▷ without a compatible software simulator
 - ▷ by only observing

- ▶ it is hard to make sure the bug really gets fixed without a test-bench
- ▶ it is slow to get bugs fixed
 - ▷ when the application is dedicated to physics and not general
 - ▷ when the engineer do not have the same physics environment to reproduce the bug
 - ▷ when the physicists can only tune limited registers to play with the bug



Linux on FPGA SoC

- ▶ SoC: system on a chip
 - ▷ CPU and system memory on a single integrated circuit
 - ▷ powerful enough to run an operating system
- ▶ features are implemented in software include OS
 - ▷ easy to learn, easy to develop and debug
 - ▷ CPU is general purpose, which runs much slower than present register-transfer level implementation
- ▶ try Linux in order to make a balance with present RTL design
- ▶ Linux is a portable open-source operating system kernel
 - ▷ there are ongoing efforts to run Linux on Virtex 2 pro and spartan 3(E) FPGAs

Results: Linux on Virtex 2 Pro

ISE	10.1
Linux	3.0
toolchain	compiled
gcc	4.5.3
binutils	2.22
glibc	2.12
rootfs	nfs
CPU	PowerPC405
clock	100Mhz
console	RS232



- ▶ user FPGA on main board
 - ▷ software level on board control
 - ▷ read data from system FPGA via board circuit and push out by Ethernet
- ▶ preliminary speed test
 - ▷ 100Mbps on board connected to 1Gbps host
 - ▷ standard FTP protocol
 - ▷ 629 KB/s, bottleneck being CPU

Results: Linux on Spartan 3E

ISE	13.2
Linux	3.0
toolchain	downloaded
gcc	4.1.2
binutils	2.16.1
glibc	2.3.6
rootfs	ramdisk
CPU	MicroBlaze
clock	50Mhz
console	RS232



- ▶ system/front-end FPGA on main board
 - ▷ replace non time critical and complex housekeeping function with software
- ▶ trigger logic/system FPGA on trigger board
 - ▷ implement smart trigger logic (e.g. FFT, DWT) with software

Conclusions and Prospects

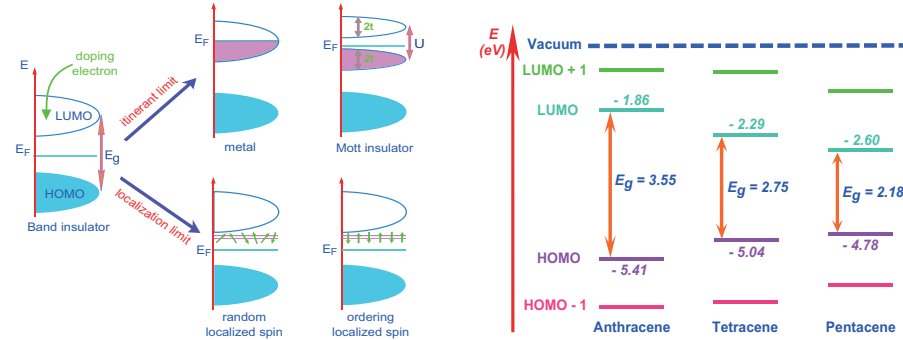
- ▶ we are able to run Linux on FPGAs on MoGURA, both hard core (powerpc) and soft core (microblaze)
- ▶ our goal of developing and debugging MoGURA firmware independent of supporting company is far from complete
 - ▷ software alone is too slow, RTL alone is too hard to master
 - ▷ somewhere in between would be the final scheme
- ▶ possibilities includes
 - ▷ hybrid: CPU architecture with software optimized by hardware accelerator wrap present RTL into accelerators
 - ▷ hybrid: independent Linux and RTL for different tasks and communicated via a bus
 - ▷ minimal change, only new features in software
 - ▷ C-to-RTL: generates VHDL from C syntax and libraries, easy to simulate e.g. No-Instruction-Set-Computer
 - ▷ myHDL: generates VHDL from python, easy to simulate
 - ▷ any mixture of the above

MAGNETIC PROPERTIES OF POTASSIUM DOPED POLYACENE: ANTHRACENE, TETRACENE, PENTACENE

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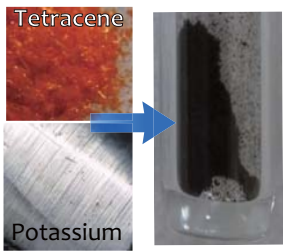
Motivation



- Searching for synthetic technique induce high homogeneity alkali metal doped organic semiconductor
- Investigating electronic states of intercalation compounds between anthracene, tetracene, pentacene and potassium

Experiment

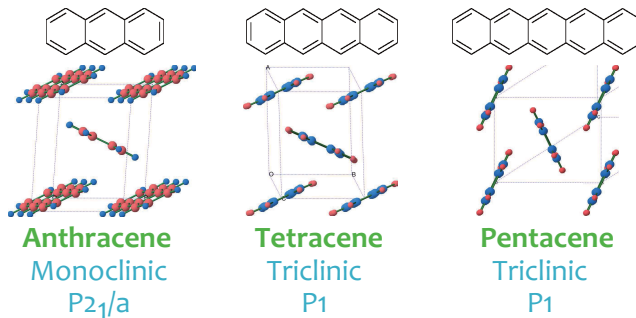
Synthesis



Annealing in glycerin bath, He atmosphere

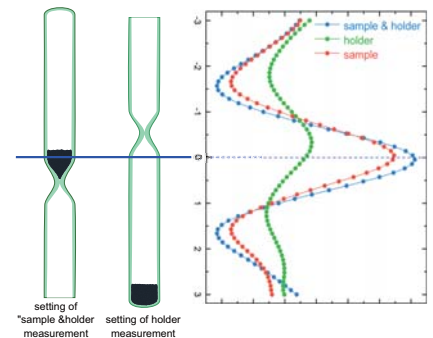
- Anthracene: 70°C
- Tetracene: 180°C
- Pentacene: 180°C

Crystal Structure & XRD measurement



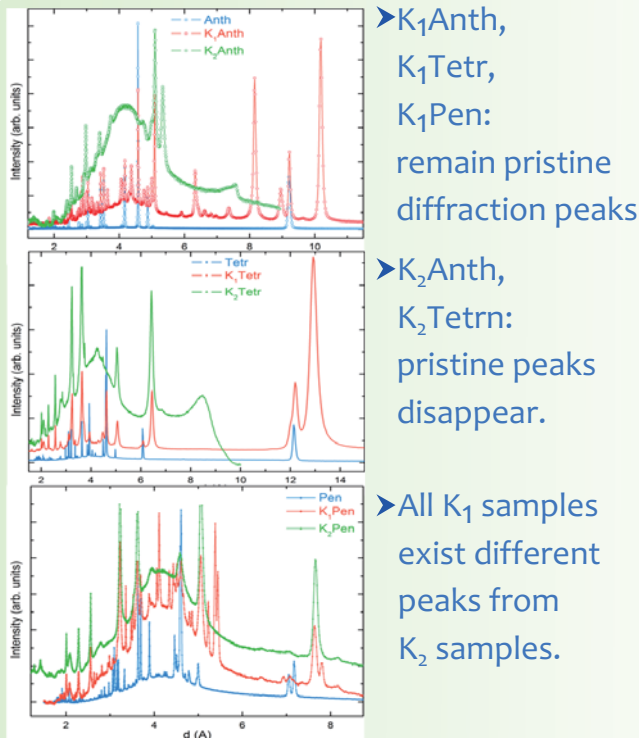
- Packed in capillary
- Condition: Ar gas
- Measured at SPring-8

SQUID measurement

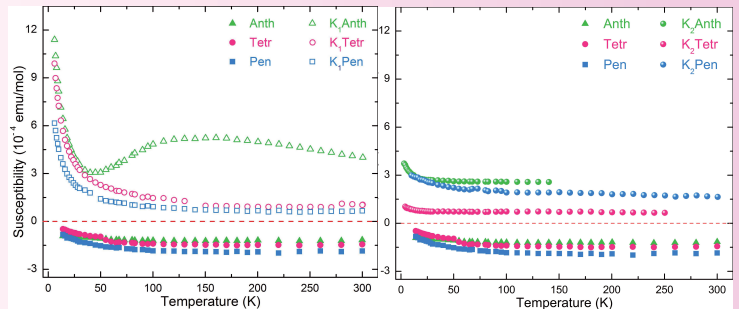


Subtract signal of quartz tube from one of sample in tube

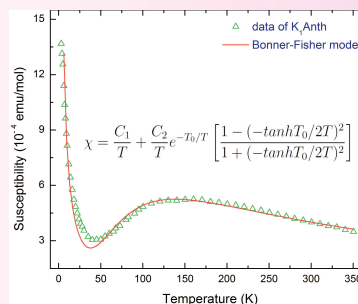
Result & Discussion



- K_1 Anth, K_1 Tetr, K_1 Pen: remain pristine diffraction peaks
- K_2 Anth, K_2 Tetr: pristine peaks disappear.
- All K_1 samples exist different peaks from K_2 samples.



- K_1 Anth show a hump K_2 Anth large magnetization
- K_1 Tetr & K_2 Tetr: same magnetization at high T range
- K_1 Pen: magnetization less than K_2 Pen



Antiferromagnetic interaction Ising chain model

Doped spin number:
 $n = 0.44$ spin/molecule
Interacting energy:
 $J = 8.3$ meV

Summary & Future work

Summary

- K doped samples with nominal composition 1 include pristine & K_2 phase beside K_1 phase.
- K_1 Anth is stable while K_1 Tetr & K_1 Pen are not.
- Antiferromagnetic interaction only exist in K_1 Anth

Future work

Synthesizing intercalation compound of other polycyclic aromatic hydrocarbons
➔ Investigating effects of molecule structure & crystal structure to electronic state

A Test for Uniformity of BGO Crystals to be Used For an EM Calorimeter

Qinghua HE, Ryuji YAMAZAKI, Shinichi MASUMOTO
Research Center for Electron Photon Science, Tohoku University

Abstract

We are constructing a 4π BGO electromagnetic (EM) calorimeter named BGOegg. It consists of 1320 BGO crystals, which are shaped as tapered blocks to make up an egg-shaped EM calorimeter. The scintillation lights of each BGO crystal will be read out through the back surface with a photomultiplier tube (PMT). If the light collection from every point along the path of interacting particles is not uniform, the whole crystal will give different responses depending not only on the energy deposit but also on the position of the light release. In other words, the amount of scintillation lights reaching PMT is affected by the shape of the BGO crystal. In general, the position dependence of light output could be attributed to the optical properties of tapered crystals and reflectors with which the crystals are wrapped. The surface treatment by an etching method can make an improvement on the uniformity. In the present test, we are measuring the effect on the uniformity due to different surface treatments and wrapping reflectors for BGO crystals. Three different surface treatments and three kind of wrapping reflectors have been tested. In this presentation, we report the results of longitudinal uniformity for each combination of surface treatment and reflector for BGO crystals.

BGOegg

The left part of figure 1 shows the overall geometry of the BGOegg covering the polar angle from 24° to 156° . This egg-shaped BGOegg consists of 1320 BGO crystals which are shaped as tapered blocks. The plot on the right side of figure 1 displays the arrangement of BGO crystals between 24° - 90° polar angle and the shape parameters for each crystal. The horizontal axis z is along the incident gamma beam as shown in the picture on the left side of figure 1 while the vertical axis represents the distance to the axis z .

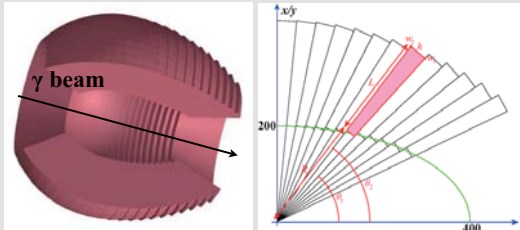


Fig. 1 The overall geometry of BGOegg (left) and the arrangement of the BGO crystals(right).

BGO Crystal

Three kinds of surface treating methods are used in order to change the optical properties. These methods are named Main, Etched and Polished respectively. More details about each method are described in the following table. Figure 2 shows the 00F BGO crystal processed with Main, Etched and Polished surface treating methods respectively. For the purpose of improving light reflection on the surfaces of the crystal, three kinds of reflectors (ESR, Al foil and Tyvek) are applied as a wrapping reflector for BGO crystal as shown in figure 3.

Surface treatment	
1.Main	Both front and back surfaces polished while the lateral surfaces etched
2.Etched	Only back surface polished while other surfaces etched
3.Polished	All surfaces polished

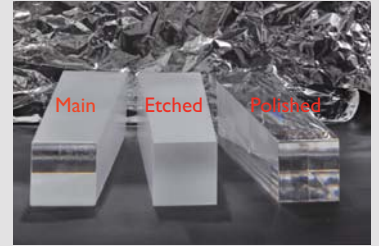


Fig. 2 Surface treatments.

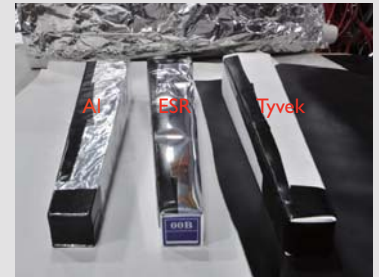


Fig. 3 wrapping reflectors.

Experimental Setup

As shown in figure 4, the experimental setup mainly consists of a photomultiplier tube, a BGO crystal, a collimator, a ^{137}Cs gamma source ($E_\gamma=661.7$ keV) and several light shielding equipments. The collimator is made of two lead blocks ($50\times100\times200$ mm) with nine holes of 6 mm in diameter, which is placed at about 120 mm far from the BGO crystal. The distance between two adjacent holes is 25 mm and the depth of the hole is 50 mm. The gamma source is placed at the entrance of one of these nine holes in order to change the measuring point on BGO crystal. A circuit diagram of DAQ system is displayed in figure 5.



Fig. 4. An photograph of the experiment setup.

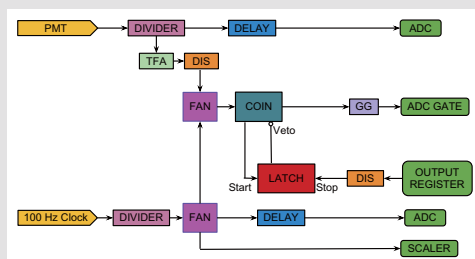


Fig. 5. Electrical data acquisition system.

Results

In this section, we report the results of the test of the effect on the uniformity of BGO crystal due to different surface treatments and wrapping reflectors. Figure 6 displays the light output as a function of the distance from PMT. Figure 7 shows the distance dependence of the resolution. The non-uniformity is plotted versus the distance in figure 8. These results show that the Main surface treatment have the best uniformity while the Polished one the worst and the ESR reflector have the largest light output, better energy resolution and an improvement on the uniformity. Based on these ideas, it can be concluded that a BGO crystal with Main surface treatment and ESR reflector has an optimal performance compared with other combinations of surface treatment and reflector.

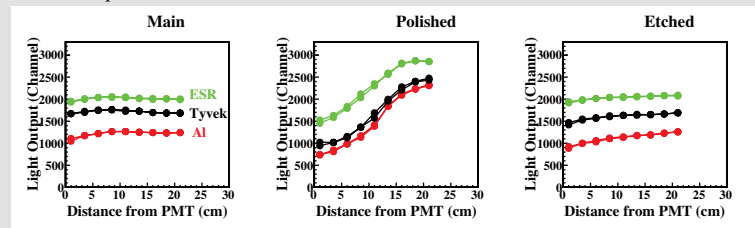


Fig. 6. Light output are plotted as a function of the distance from PMT.

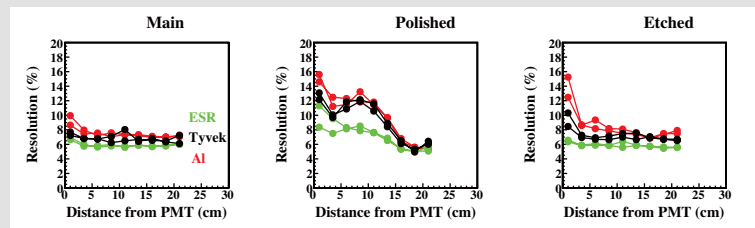


Fig. 7. Resolution is plotted versus the distance from PMT.

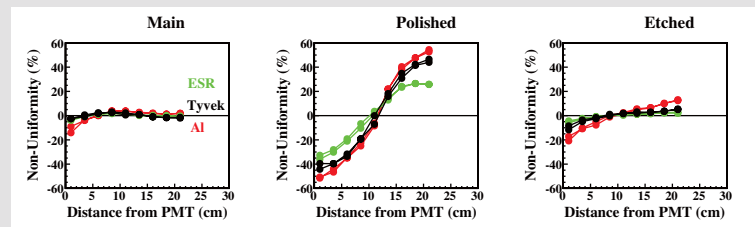


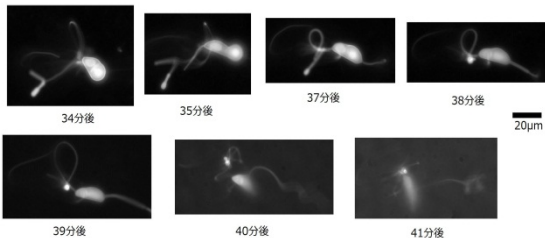
Fig. 8. Non-uniformity is plotted as a function of the distance from PMT.

A living body system uses an amphipathic molecule etc. as a fundamental component

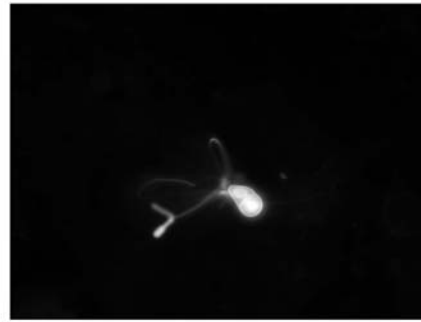
S.FUKAYA
Department of Physics, Tohoku univ., Sendai,
Japan

Time variation of vesicle

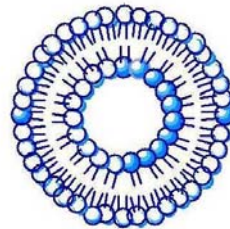
- NaOH solution was dropped



Vesicles has 'strange' form



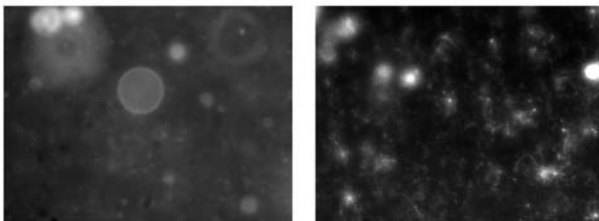
What is a Vesicle



- Spherical lipid bilayer membrane

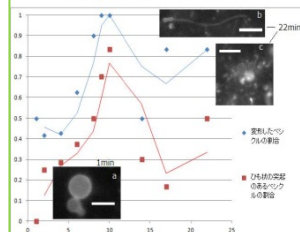
effect of changes in pH

- normal
- NaOH dropped

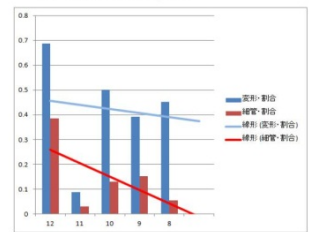


effect of changes in pH(graph)

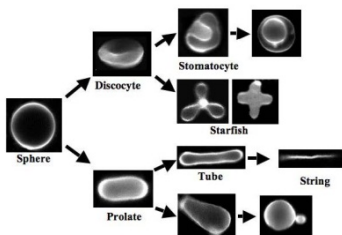
- In pH12



- In various pH



How does form of a vesicle change?



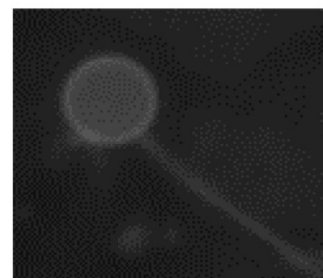
- Area
- Difference
- Elasticity
- model

$$F = F_b + F_l + F_A$$

$$F_b = \int \left(\frac{1}{2} \kappa (2H)^2 + \bar{\kappa} K \right) dA$$

$$F_l = \sum \sigma \oint dl^{(i)}$$

$$F_A = \frac{\langle \kappa_r \rangle}{2Ad^2} (\Delta A - \Delta A_0)^2$$



...but
Tubular formation has not been explained completely!

CRYSTAL GROWTH OF NEW TARGET SYSTEMS FOR HIGH-ENERGY NEUTRON-SCATTERING MEASUREMENTS AT J-PARC

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¹Department of Physics, Tohoku University, Sendai, Japan

²Institute for Materials Research, Tohoku University, Sendai, Japan

³World Premier International Research Center, Tohoku University, Sendai, Japan

Introduction

Magnetism on doped CuO₂ plane is considered to play a key role for the emergency of high-T_c superconductivity. Therefore, the observation of spin fluctuations is important. However, huge amount of crystal is required for the neutron scattering measurement due to the low flux.



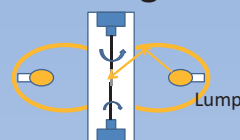
~10cc

Motivation

Except for few examples, growth of large single crystal of high-T_c cuprate is difficult. To overcome difficulties in the crystal growth, and to obtain sufficiently large crystal for the neutron scattering measurement, we have attempted improved floating-zone method for a couple of systems. Also, we have prepared high-quality but tiny crystals and assembled them to gain the effective sample volume.

Method

Floating-Zone method

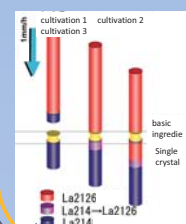


Merit

Large and impurity-free crystal can be available.

Demerit

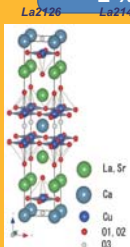
Difficulty in finding out the growth condition



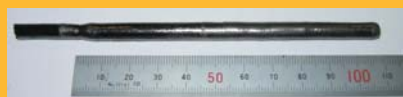
We used a feed rod with a concentration-gradient. The crystal growth was started from a part for which the growth is rather easy and continuously proceeded with changing conditions.

Crystal Growth

La_{2-x}Ca_{1+x}Cu₂O₆

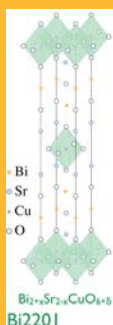


basic ingredient	La2126-System	La214-System
Solvent composition	La _{1.88} Ca _{0.12} Cu ₂ O ₆	La _{1.88} Ca _{0.12} Cu ₂ O _{6+δ}
Solvent quantity	200-270mg	200mg
Air	Oxygen flow	Oxygen flow
Feed	1.00mm/h	1.00mm/h
Rotation	25rpm	10rpm
Lump power	160-170W × 2	200W × 2



Large single crystals for neutron scattering measurement have been successfully obtained.

Bi_{2+x}Sr_{2-x}CuO_{6+δ}

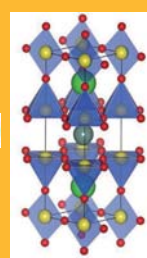


basic ingredient	Bi2201-System
Solvent composition	None
Air	Oxygen flow
Feed	0.60mm/h
Rotation	
Lump power	150W × 2



YBa₂Cu₃O_{6+δ}

YBa₂Cu₃O₇



3.0 × 3.0 × 0.3[mm]

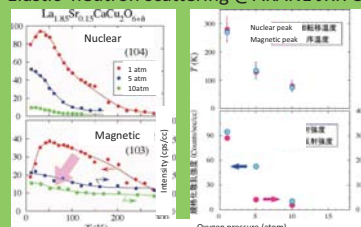


10pieces ~ 100mg

Several tiny crystals were assembled.

Results of Preliminary Neutron Scattering measurement

Elastic neutron scattering @AKANE JRR-3



With increasing the oxygen pressure during the crystal growth, intensity of both nuclear and magnetic peaks are reduced.

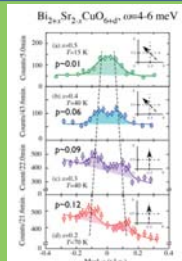
Close relation between crystal structure and magnetic order



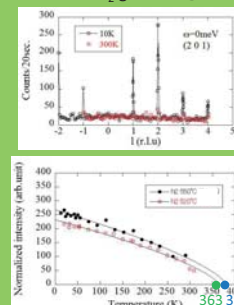
The first observation of spin fluctuation in Bi2201

Incommensurate spin fluctuations exist and the peak-splitting increases upon doping, similar to the case of LSCO

Possible universality of spin correlation in single-layer system



annealed in N₂ gas flow @530° C



Magnetic order and its doping-dependence were confirmed

Summary

Utilizing an improved growth technique, we have succeeded in growing large single crystal of new target systems for neutron scattering measurement. Preliminary measurements have been done at triple-axis spectrometer and we observed magnetic response in the new samples. Based on these results, we are going to study high-energy spin dynamics by using a newly constructed high-flux neutron beam facility, J-PARC



Simple model for rupture process of pressure-sensitive adhesives

Shinobu Sekine

Tohoku University, Department of physics

2012/02/20

Rupture block model

Our model is extended from Yamaguchi's model¹. Original Yamaguchi's model is not treated in the final rupture process.

Major features of this model:

- 1 the adhesive layer is divided into N_b rectangular blocks of equal size
- 2 the flow of films is assumed to be superposition of slippage and parabolic deformation of the block
- 3 cavity is assumed to be described by the gap between adjacent blocks
- 4 cavitation dynamics is described by the Rayleigh-Plesset equation
- 5 If the elongation ratio is over the limitation condition, this block is divided into half.

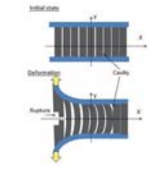


Figure: Schematics of PSA rupture block model

¹T.Yamaguchi et al:Eur.Phys.J.E 20,7-17(2006)

Surface tension and rupture process.

Surface tension is acting on upper side of blocks because of the deformation.

$$\eta \dot{z}_j = \gamma \frac{(z_{j+1} + z_{j-1} - 2z_j)}{w_0^2}, \quad z_j(t) = H_0 \lambda_j(t)$$

$$\therefore \lambda_j = \frac{\gamma}{\eta w_0^2} (\lambda_{j+1} + \lambda_{j-1} - 2\lambda_j) \quad (2)$$

$$\eta : \text{viscosity}, \gamma : \text{surface energy} \quad (3)$$

R.H.S. is assumption of curvature elasticity.

To describe rupture process, we introduce the forming rupture condition:

$$\lambda(t)_j > \lambda_{\text{limit}} = 2.0 \quad (4)$$

If one of these block's elongation ratio is over this limitation, this block is ruptured at the center point. After this rupture process, the under surface of this block, which is created this process newly, has same motion under surface tension describe above.

Effect of cavity

To consider the effect of cavity, we compare CASE1, which is cavity existed, and CASE2, no cavity existed.

case	cavity
CASE1	exist
CASE2	not exist

Cavity has very important effect for rupture process to weak the stress of the adhesive layer.

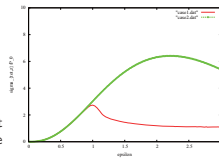


Figure: The stress-strain (S-S) curves for the effect of cavity

Effect of separation speed(1/2)

Deborah number

$$D_e = \frac{\tau \dot{\lambda} L_0}{H_0} = \frac{\text{stress relaxation time}}{\text{time scale of observation}} \quad (5)$$

At higher D_e , fluid property changes viscosity to elasticity.

The peak stress value increases as the separation speed $\dot{\lambda}$ goes up.

case	$\dot{\lambda}$	D_e
CASE1	0.1	500
CASE5	0.5	2500
CASE6	0.02	100

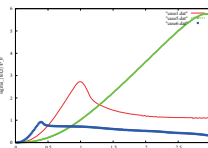


Figure: S-S curves for various $\dot{\lambda}$

Conclusion

- We extended Yamaguchi's block model to treat rupture process.
- In qualitative view point, our results are good agreement of experimental data.
- Cavity has very important effect for rupture process to weak the stress of the adhesive layer. As can be seen from this result, to compare simulation results and experiments quantitatively, we would need more detailed modeling, for examples, shape and position of cavity, stress concentration around cavity.
- The slip velocity is proportional to the shear stress at the interface. $\mu \dot{\lambda}_{sl} = \sigma_{sl}$. Though such model is standard for simple liquid, non-linearity and memory effect will be important for adhesives.
- The constitutive model is the viscoelastic fluid one, while the commercial PSA is viscoelastic solid because it is actually weakly cross-linked.

Abstract:

Pressure-sensitive adhesives (PSA) are very useful in our ordinary life, which is a thin layer applied, for example, on the surface of tape film. Typical materials of PSA are made of block copolymers such as acrylates or styrene-isoprene-styrene (SIS) triblocks. They are usually very soft and highly dissipative, and can stick on a variety of surfaces under low pressure in short time without any solvent evaporation, heating process or chemical reactions.

We propose a simple mechanical model describing viscoelasticity and cavitation during the crack propagation process in pressure-sensitive adhesive. This model is originally proposed by Yamaguchi et al., They applied this model for probe-tack test, but their model is not applicable to the final rupture process. We extended this model to treat this rupture process, and applied to the situation of peel test.

In this presentation, we will report our extended model and calculation results.

Deformation of one block

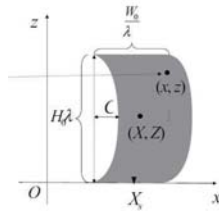


Figure: Coordinate system to describe the block motion and deformation

- (X, Z) : Coordinates of the center of gravity,
- X_c : Coordinate of the central position of the side surface.
- C : Parameter characterizing the parabola \propto curvature of the side surface
- H_0 : Initial thickness
- W_0 : Initial width of one block
- $\lambda(t)$: Elongation ratio

Parameter

The typical set of parameters.

Parameter	Unit	Value	Comment ³
PSA thickness H_0	m	1.0×10^{-4}	
PSA width L_0	m	5.0×10^{-3}	
Separation speed V_z	m/s	1.0×10^{-4}	$\dot{\lambda} = 0.1$
Atmospheric pressure P_0	Pa	1.0×10^5	
Viscosity η	Pa s	2.0×10^5	water 1.0×10^{-3}
Elastic modulus G	Pa	2.0×10^3	Au 2.7×10^{10}
PSA surface energy γ	J/m ²	30.0×10^{-3}	water 70.0×10^{-3}
Friction coefficient μ	Pa s/m	2.0×10^9	
Initial cavity radius R_0	m	1.0×10^{-6}	
Number of blocks N_{block}		100	
Rupture condition λ_{limit}		2.0	

$$\text{Relaxation time } \tau = \frac{\eta}{G} = 100[\text{s}]$$

³material properties are referenced by Rika Nenpyo

Effect of thickness(1/2)

case	$H_0[\text{m}]$	H_0/L_0
CASE1	1.0×10^{-4}	0.02
CASE3	2.0×10^{-4}	0.04
CASE4	5.0×10^{-5}	0.01

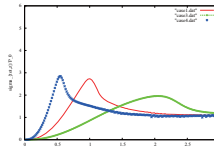


Figure: The S-S curves for various thickness

- 1 Initial stage is strongly affected by the thickness of the adhesives.
- 2 The stress at the peak decreases as the thickness of the adhesives increases.
- 3 The peak position is dependent of the film thickness.

Effect of separation speed(2/2)

Origin of these results

- Higher separation rate causes larger shear stress and larger resistance for the cavity expansion, leading to the larger negative pressure (6).

$$\dot{R}_j = \frac{R_j}{2\eta} (P_{\text{cav},j} - P_j) - \frac{\gamma}{2\eta}$$

$$\therefore P_j = P_{\text{cav},j} - \left(\frac{2\eta \dot{R}_j}{R_j} + \frac{\gamma}{R_j} \right) \quad (6)$$

Lead to large negative pressure

PSA debonding process

PSA=Pressure Sensitive Adhesives.

Example of PSA is the viscoelastic surface layer of packing tape.

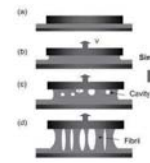


Figure: Schematics of PSA debonding process

The debonding process of PSA²:

- Initial state
- Uniform deformation
- Cavity expansion
- fibrillation
- fracture of fibrils or debonding from the probe

²T.Yamaguchi et al:Eur.Phys.J.E 20,7-17(2006)

Constitutive equation

The upper-convected Maxwell fluid model

$$\frac{\partial \sigma_{ij}}{\partial t} = -\frac{1}{\tau} \sigma_{ij} + \left(\frac{\partial v_k}{\partial x_i} \right)^T \cdot \sigma_{kj} + \sigma_{ik} \cdot \left(\frac{\partial v_k}{\partial x_j} \right) + \left(\frac{\partial v_j}{\partial x_i} + \frac{\partial v_i}{\partial x_j} \right)^T G \quad (1)$$

- With a single relaxation time $\tau = \frac{\eta}{G}$
- Generalization of Newtonian fluid to satisfy principle of material objectivity.

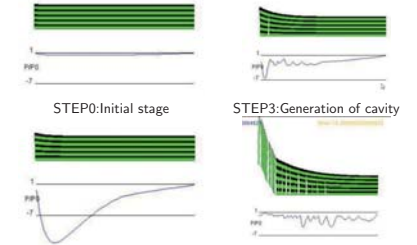
Principle of material objectivity

Physical properties, such as angle, length, are unchanged by any parallel displacement and rotation as a whole in any manner in space²

²T.Kyoja: Continuum Mechanics (in Japanese), Morikita, 2008

Snap shots of deformation

For the symmetry of the model system, these images are drawn on



STEP2: Early state of deformation STEP4: Beginning of Rupture

Effect of thickness(2/2)

Origin of these results

From the detail calculation of upper convected Maxwell model (1), we get the following result:

$$\frac{\partial \sigma_{s,j}}{\partial t} = -\frac{\sigma_{s,j}}{\tau} + \frac{4\dot{C}_j}{H_0 \lambda} (\sigma_{zz,j} + G)$$

In the initial stage, $\sigma_{s,i} = 0$:

$$\frac{\partial \sigma_{s,j}}{\partial t} \approx \frac{4\dot{C}_j}{H_0 \lambda} (\sigma_{zz,j} + G)$$

where σ_s : shear stress at the surface.

R.H.S shows that thinner the thickness of the adhesives H_0 , larger the rate of stress increases.

Effect of friction coefficient

The stress decreases with the decrease of μ since slip-page at the surface reduces the negative pressure.

case	μ
CASE1	2.0×10^9
CASE7	2.0×10^8
CASE8	2.0×10^{10}

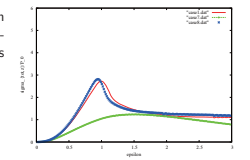


Figure: S-S curves for various friction coefficients

An investigation of the photo induced production of strangeness in the threshold energy range.

B. Beckford for NKS2 collaboration, Dept. of Physics, Tohoku University

Introduction

Strangeness photoproduction

Strangeness production on a nucleon or a nucleus by the electromagnetic interaction provides invaluable information on the strangeness production mechanism and strengths of meson-hadron coupling constant. Kaon photoproduction can be a good probe to find missing resonance states.

NKS2 experiment aims to explore the strangeness photoproduction on the deuteron by measuring neutral kaons and Λ hyperons with tagged photon beams of $E_\gamma = 0.8$ -1.1 GeV.

$n(\gamma, \Lambda) K^0$ reaction in the threshold region

$n(\gamma, K^0)\Lambda$ process has unique features in the investigation of kaon production process by electromagnetic interaction as follows,

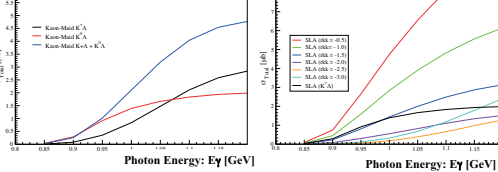
- no charge in initial and final state
 - \rightarrow t-channel Born term does not contribute
- isospin symmetry to $p(\gamma, K^0)\Lambda$ process
 - \rightarrow sign of coupling constant in u-channel is opposite
- the electromagnetic coupling constants of resonances in the s- and t-channels
 - \rightarrow h^+ s different from K^+ process : e.g., $g(N^*\eta)$ and $g(K^*\Sigma^0 p)$

Due to these characteristics, the interference among the diagrams in the K^0 production process is quite different from that in the K^+ process. For the elementary reaction of $\gamma n \rightarrow K^0 \Lambda$, photon energy dependence and angular distribution at $E_\gamma = 1.05$ GeV are calculated using Kaon-MAID model and Saclay Lyon model (SLA). These two isobar models agree well in their predictions of the $\gamma p \rightarrow K^+ \Lambda$ process, however they are significantly different for the $\gamma n \rightarrow K^0 \Lambda$ process.

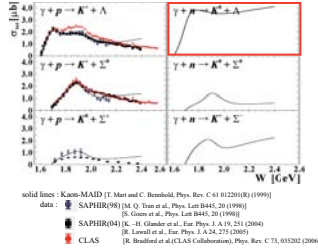
Total Cross Section Incident Photon Energy Dependence

Kaon-MAID and SLA Isobar Models

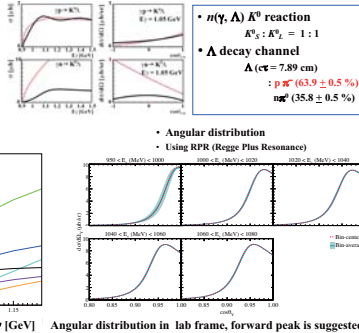
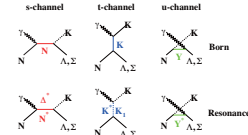
Total Cross Section of Inclusive A Measurement



Total cross section for 6 isospin channels of kaon photo-production



Feynman Diagrams for isobar model



Experimental Apparatus

Tagged photon beam

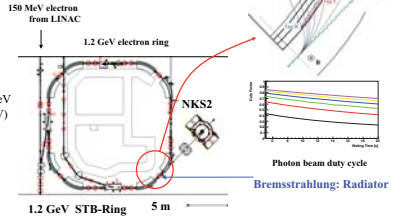
NKS2 was installed at Laboratory of Nuclear Science, Tohoku (LNS-Tohoku).

The photon beam is generated via bremsstrahlung and the scattered electron's energy is tagged by STB-Tagger system.

- Energy Region 0.8 - 1.1 GeV
- $n(\gamma, K^0)\Lambda$ threshold (915 MeV)
- Energy Resolution 10 MeV
- Beam Intensity 2×10^6 Hz
- Duty Factor 35 - 85 %

Schematic view of the experimental hall in ELPH

STB Tagger System

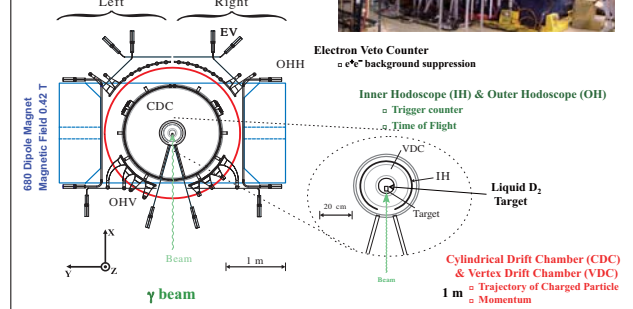


NKS2 Upgrade Setup

Λ are measured with Neutral Kaon Spectrometer (NKS2) by detecting the $\Lambda \rightarrow p \pi^-$ decay channel.

NKS2 is a spectrometer dedicated to measure the $n(\gamma, \Lambda) K^0$ reaction in the threshold region. It consists of a dipole magnet and two types of drift chamber, TOF counter. The liquid deuterium target is located at the center of NKS2.

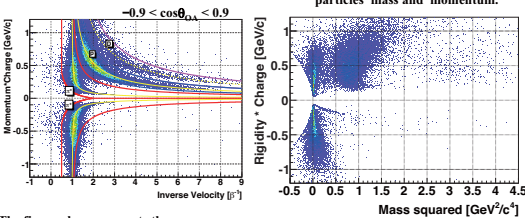
Schematic view of NKS2



$p\pi^-$ Event Distribution

Particle Identification (PID)

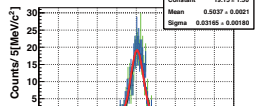
PID capability of two charged particle track events



The correspondence particles between the square of the particles' mass and momentum.

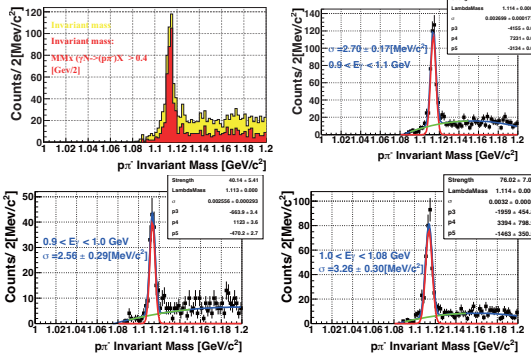
The figure above presents the correspondence particles between momentum and inverse velocity. The red, blue and magenta line represents the pion, proton and deuteron designated regions.

The figure at the right is the kaon missing mass in the MMx ($\gamma N \rightarrow (p\pi^-)X$) measurement



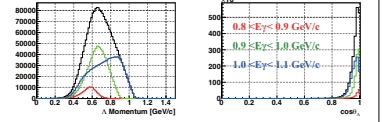
Invariant Mass Distribution

This figure presents the invariant mass distribution of $p\pi^-$ each energy bins. The $p\pi^-$ invariant mass is shown with the missing mass region of $MM_{\text{miss}} > 0.40$ GeV/c^2. The green line shown in the fit represents the estimated background.



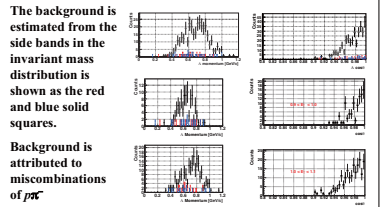
Momentum and Angular Distribution

Simulated momentum and angular distributions



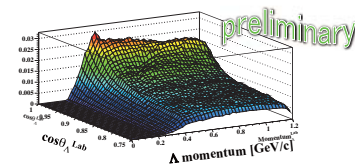
The Lambda particle was generated isotropically in the laboratory frame. The distribution is shown for the energy range of 0.8-1.0, 0.8 - 0.9, 0.9-1.0, 1.0-1.1 GeV as black, red, green and blue lines respectively.

Measured momentum and angular distributions



Detector Acceptance

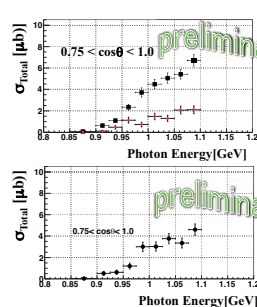
Λ Acceptance



Simulation Condition :

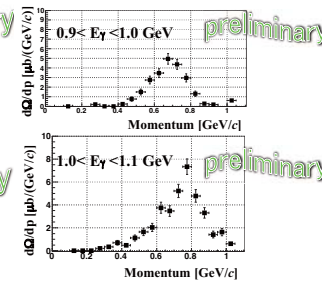
Generated Λ isotropically in Lab frame,
 $0 < p_\Lambda < 1.2$ GeV, $0.75 < \cos \theta_{\Lambda} < 1.0$
Use the same analyzer for the experimental data.

Total Cross Section



Momentum Spectra

After acceptance correction



Summary

- The $\gamma + n \rightarrow K^0 + \Lambda$ process plays a unique role in the investigation of photoproduction mechanism
- We have successfully the upgrade of Neutral Kaon Spectrometer (NKS2) The experiment in the photoproduction of strangeness on a deuteron target was accomplished using a tagged photon beam at the ELPH research facility
- The invariant mass spectrum of $p\pi^-$ was obtained and the resolution of the spectrums have been detailed
- The momentum distribution was obtained for two photon energy region, 0.9 to 1.0 GeV and 1.0 to 1.1 GeV.
- Preliminary results have been reported.

Research of Transmission-line STJ Detector for Terahertz Band

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^a RIKEN, ^b Tohoku University, ^c National Astronomical Observatory of Japan, ^d Techno X Co.,Ltd.



東北大学

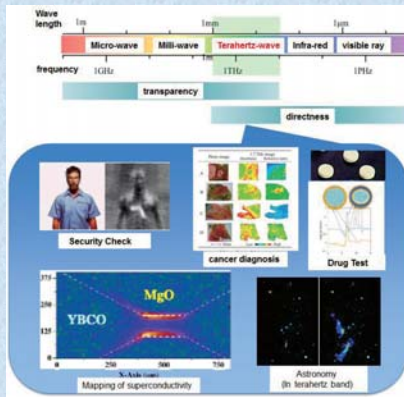


Abstract

We have proposed and demonstrated a new broadband and high efficiency THz wave detector using a superconducting tunnel junction (STJ). The detector uses by Cooper-pair Braking (CPB) process. The detector consists of two log-periodic antenna wings connected with an impedance transformer and two long (micro-stripline) STJs. The detector utilizes the excess tunneling current caused by the Cooper pair breaking due to the radiation whose photon energy is greater than the energy gap of the used superconductor (around 0.7 THz for niobium). We have fabricated the device and performed the experiments for confirming the principle of this detector. We have succeeded the detection for the first time.

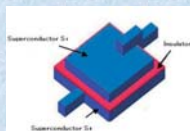
0. Introduction

Terahertz (THz) waves are electromagnetic waves with frequencies between high-frequency edge of the microwave band and the long-wavelength edge of far-infrared light. THz waves were non-development from difficulty of the generation and the detection, but a study of THz band advances in late years, THz band attracts attention in the field of industry and research.

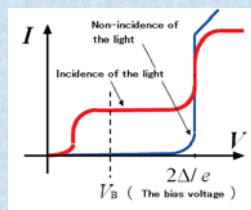
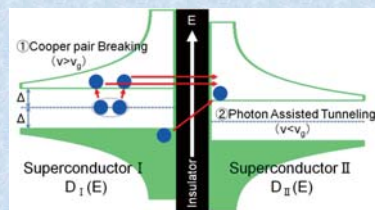


1. Design of Micro-stripline Detectors

Mechanism of STJ detector



There are two detection processes to the STJ detector. One is the Cooper-pair Braking (CPB). The other is the Photon Assisted Tunneling (PAT) process. They are separate by the gap frequency of the superconductor.

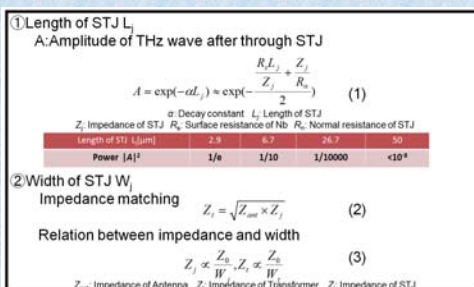
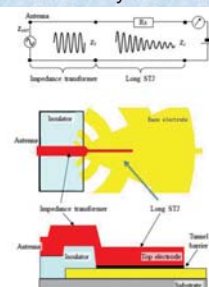


Density of the state of the STJ

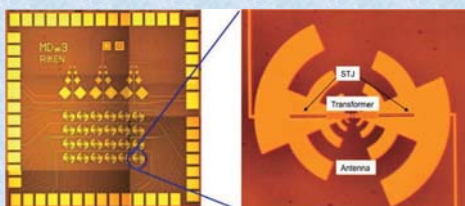
I-V curve of the STJ

Transmission-line STJ Detector

If the length of STJ is enough long compared with the mean free path of the emergent quasiparticles (~10 μm for niobium), the energy of THz waves is absorbed by the STJ behaves as a transmission line and we detect it.



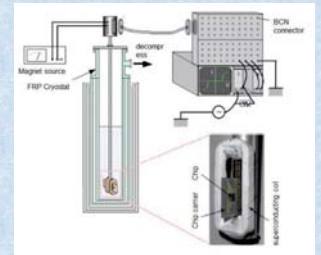
2. Fabrication of STJ



We have fabricated the STJ detector in RIKEN clean room. The left photograph is the general view and magnification image.

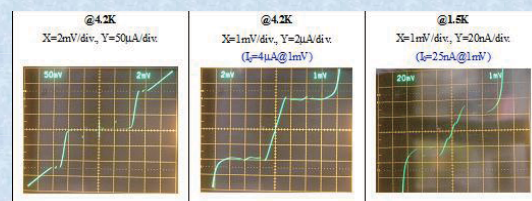
L_j : 4 kinds \times W_j : 3 kinds = 12 kinds of combinations of length and width

3. Current-Voltage measurements



We measure the performance of the element by current voltage measurement. An important point is a value of the subgap leak currents of 4.2K and 1.5K. Because the decline of the leak current means the decline of the noise.

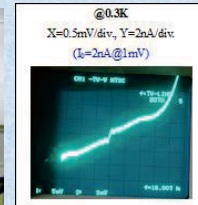
4 μm × 27 μm-STJ detector



This STJ has good I-V curve in the subgap leak current. The current of 1.5K becomes one hundredth in comparison with the current of 4.2K.

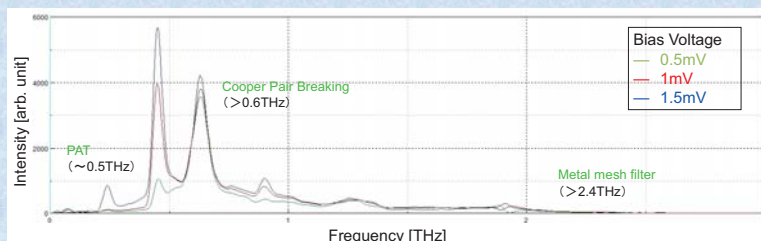
4. Optical measurements (FT-IR)

Fourier transform infrared spectrometer and 0.3K cryostat



We measured optical properties with Fourier transform infrared spectrometer. We measured in 0.3K so that the leak current reduce enough.

4 μm × 27 μm-STJ detector



- CBP process detection is seen in more than 0.6THz.
- As bias voltage increases, PAT process detection increases.
- More than 0.7 THz detection is disturbed by the niobium antenna.

5. Summary

1. We have proposed a new THz wave detector using two long micro-stripline STJs.
2. We fabricated the device and performed the experiments for confirming its detection principle.
3. We have succeeded the detection of THz wave for the first time.

References:

1. T. Noguchi, T. Suzuki, 10th Workshop on Submillimeter-Wave Receiver Technologies in Eastern Asia (2009).
2. K. Takahashi, Master Thesis (in Japanese) (2011).
3. S. Ariyoshi et al., EUCAS2011 (2011).

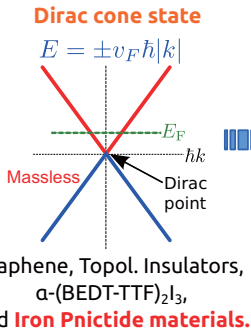
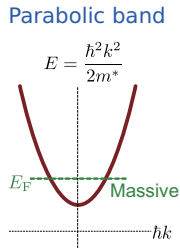
Evidence for Quantum Magnetotransport of Dirac Cone States in Ba(FeAs)₂

K. K. Huynh^{a,*}, Y. Tanabe^{b,*}, T. Urata^a, R. Nouchi^b, N. Mitoma^a, S. Heguri^a, J. Xu^b, G. Mu^a, H. Nohjiri^c, K. Tanigaki^{b,*}

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Introductions and motivations

A Dirac Cone State (DCS) is a quantum state with a linear dispersion relation between energy (E) and momentum (k).



1 Landau level (LL) splitting in magnetic field B:

$$E_n = \pm v_F \sqrt{2\hbar e B |n|}$$

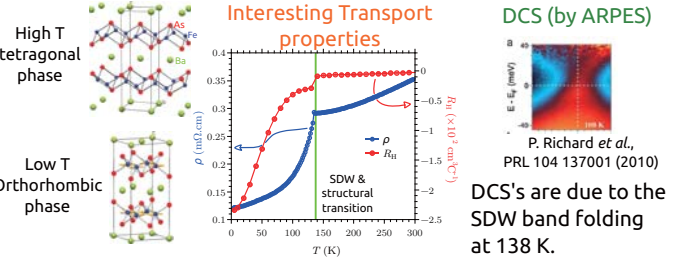
Distinguished energy scale → large LL splitting in a small B.

Intriguing quantum magnetotransport phenomena (e.g. Quantum Hall Effects at room temperature)

2 Very light mass: high mobility carriers.

* Ba(FeAs)₂ material

DCS is theoretically predicted and experimentally observed (ARPES) in Ba(FeAs)₂. In this material, the DCS's appear as SDW gap nodes.



- 1 Effects of DCS's on transport properties of Ba(FeAs)₂?
- 2 Possible quantum transport of DCS's?

Experimental Results

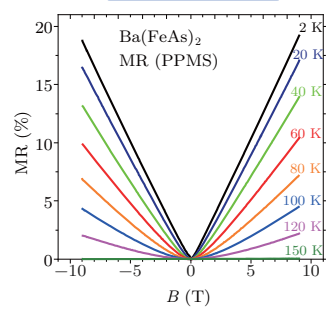
Transport properties, including Resistivity (ρ), Magnetoresistance (MR), and Hall coefficient (R_H) of Ba(FeAs)₂ single crystal were measured and analyzed in order to understand of the transport properties of the DCS's.

* Experimental setup

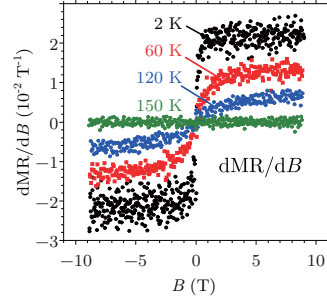
Single crystal grown out from FeAs flux.
ρ, MR, and R_H were measured using PPMS.
High-B (B ≤ 17.5 T) MR were measured at High Field LAB, Tohoku University.

* Interesting magnetotransport property

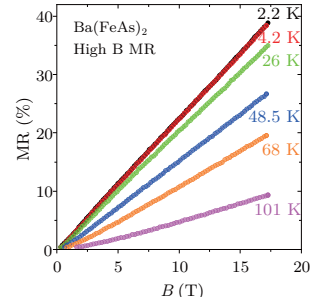
MR is linear.



dMR/dB demonstrate clearly the B-linear dependence.



Linear MR persists up to B ≤ 17.5 T



For an usual semimetal, MR is B²-dependent. The B-linear MR indicates a distinct and new transport phenomenon.

Analyses and Discussions

A detailed investigation of the B-dependent MR and dMR/dB reveals that the magnetotransport is composed of two regimes: **the semiclassical one at low B** and **the quantum regime at higher B**.

Low B-region: Semiclassical transport *

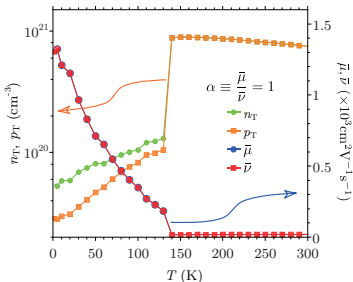
Two-carrier-type semiclassical model:

$$\rho(0) = \frac{1}{e(n_T \bar{\mu} + p_T \bar{\nu})}$$

$$MR = \frac{\rho(B) - \rho(0)}{\rho(0)} = \frac{n_T p_T \bar{\mu} \bar{\nu} (\bar{\mu} + \bar{\nu})^2}{(n_T \bar{\mu} + p_T \bar{\nu})^2} B^2$$

$$R_H = \frac{-n_T \bar{\mu}^2 + p_T \bar{\nu}^2}{e(n_T \bar{\mu} + p_T \bar{\nu})^2}$$

Total numbers (n_T & p_T) and effective mobilities (μ̄ & ν̄) of electron and hole can be estimated.



Carrier numbers drop down by 90% whereas mobilities of both electrons and holes jump up by 10 times at the SDW transition.

Both electron-like and hole-like DCS's.

Two transport regimes

* High B-region: Quantum transport

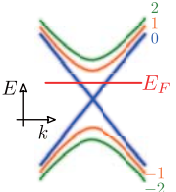
Splitting between the 0th and the 1st LL's:

$$\Delta_1 = |E_{\pm 1} - E_0| = \pm v_F \sqrt{2\hbar e B}$$

A. A. Abrikosov, PRB 58 2788 (1998)

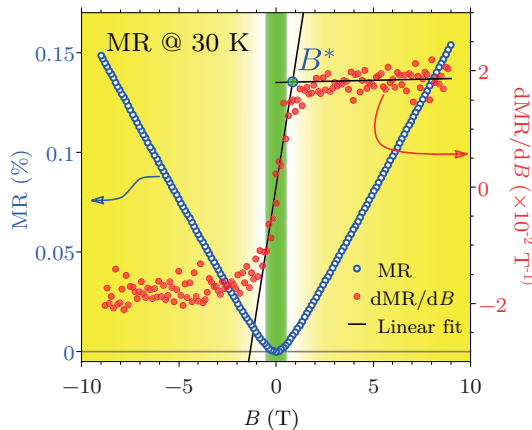
When Δ₁ is larger than both Fermi and thermal energies, only the 0th LL is occupied.

$$\begin{cases} \Delta_1 > E_F \\ \Delta_1 > k_B T \end{cases}$$



$$MR \propto (N_i / e n_D^2) \times B \quad (N_i: \text{impurity})$$

Linear MR is the transport at the limit of 0th LL in DCS's.



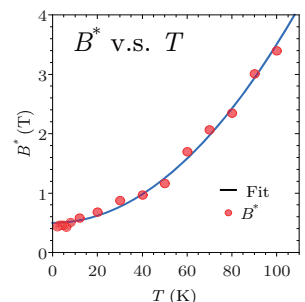
* The semiclassical transport regime (MR ∝ B²) crossovers to the quantum MR regime (MR ∝ B) at B*.

The 1st LL splitting equals to E_F and k_BT at B*:

$$B^* = (1/2e\hbar v_F^2) (k_B T + E_F)^2$$

Estimations of E_F and v_F

$$E_F = 1 \pm 5 \text{ meV} \quad v_F \approx 1.88 \times 10^5 \text{ ms}^{-1}$$



Conclusions

- 1 Below the SDW transition, transport properties of Ba(FeAs)₂ are dominated by the DCS's.
- 2 The coexistence of electron-like and hole-like DCS's were observed. These two DCS's may reside at different k_z positions.
- 3 Quantum transport of the 0th LL in the DCS's are observed and confirmed by using the linear MR and the Abrikosov's model.

Double Chooz: A Search for the Neutrino Mixing Angle θ_{13}



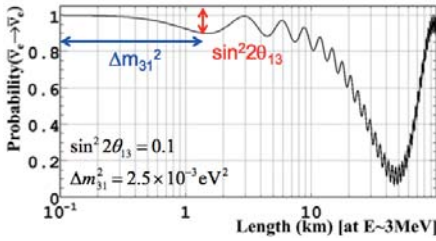
Thiago J. C. Bezerra
on Behalf of the Double Chooz Colaboration
thiago@awa.tohoku.ac.jp
Research Center for Neutrino Science
Tohoku University



Purpose and Motivation

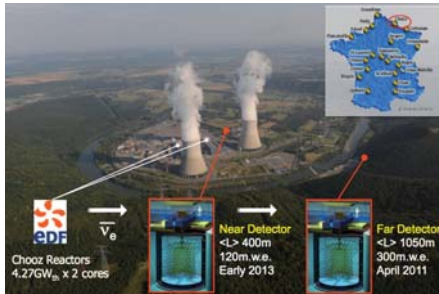
The main purpose of the Double Chooz Experiment is to determine the last unknown Neutrino Mixing Angle, θ_{13} . It is important to know this parameter since it will let to the determination of the δ_{CP} in the neutrino sector and mass hierarchy.

$$P(\bar{\nu}_e \rightarrow \bar{\nu}_e) \cong 1 - \sin^2 2\theta_{13} \sin^2 \left(\frac{\Delta m_{31}^2 L}{4E} \right) + O(10^{-3})$$



Method

To obtain the θ_{13} value, the experiment will use two identical detectors, in order to cancel the systematics uncertainties (neutrino flux, interaction cross-section, number of target and efficiency).



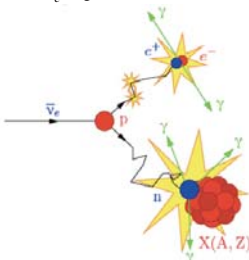
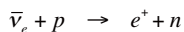
Expected Antineutrino Rate and Detection Principle

Electron antineutrinos are produced in nuclear reactors through beta decays from the fission fragments. The two Chooz commercial reactors are used as source of antineutrinos and the expected number of events can be predicted using the following relation:

$$N_v^{\text{exp}}(E, t) = \frac{N_p \epsilon}{4\pi L^2} \times \frac{P_{th}(t)}{\langle E_f \rangle} \times \langle \sigma_f \rangle$$

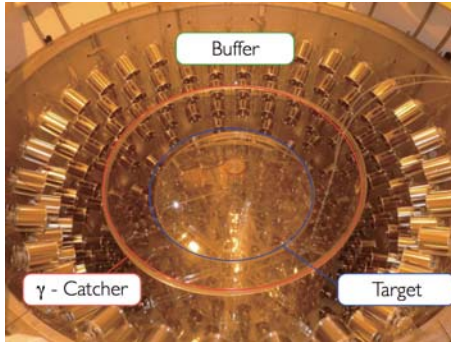
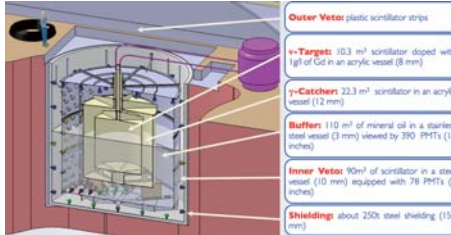
where N_p is the number of protons in the detector, $\langle E_p \rangle$ is the mean energy released per fission and $\langle \sigma_f \rangle$ is the mean cross section per fission.

At the detector, these antineutrinos will interact with the protons of the liquid scintillator, through the reaction:



The positron annihilation and the neutron capture give the signature of an antineutrino interaction.

Detector Design



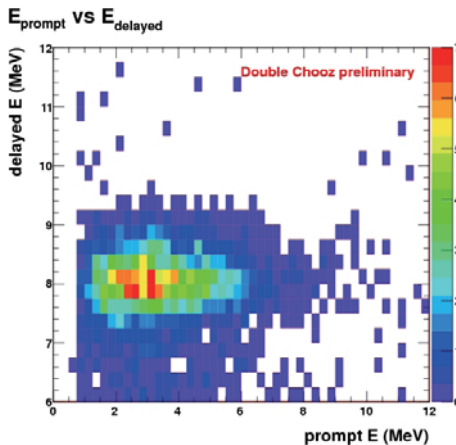
Data Taking and Event Selection

Since mid-April of 2011, the Double Chooz Far Detector is operational and getting data continuously. Its average total efficiency for physics runs, i.e. not considering calibration periods, is 77.5% of the time. For the first result, a total run time of 101.523 days were used, corresponding to a live time of 96.823 days, due to 1 ms of muon veto.

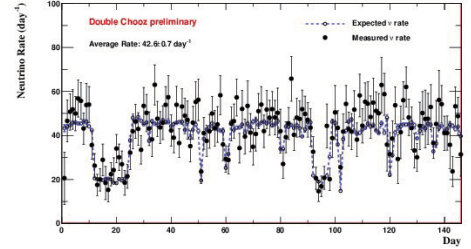
The main sources of background can be suppressed by a delayed coincidence selection:

- Prompt signal:** positron energy ($0.7 < E_p < 12$ MeV)
- Delayed signal:** gamma emission when a neutron is absorbed by the gadolinium ($6 < E_D < 12$ MeV)
- Time correlation:** $\tau \sim 30 \mu\text{s}$ ($2 < \tau < 100 \mu\text{s}$)

In addition, some qualities cuts are also applied.



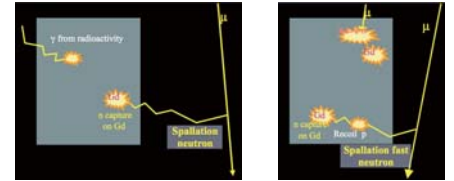
Neutrino candidates rate (background not subtracted)



Backgrounds

The backgrounds can be divided in two types:

- Accidental: environmental gamma-ray (prompt like) and muon induced neutrons (delayed). It can be estimated by off-time method.
- Correlated: caused by fast n or $^9\text{Li}/^8\text{He}$. Fast n will recoil a proton (prompt like) and, after thermalized, will be captured on Gd. $^9\text{Li}/^8\text{He}$ are beta-n emitters, produced by muon-induced spallation. It cannot veto due to long lifetime.

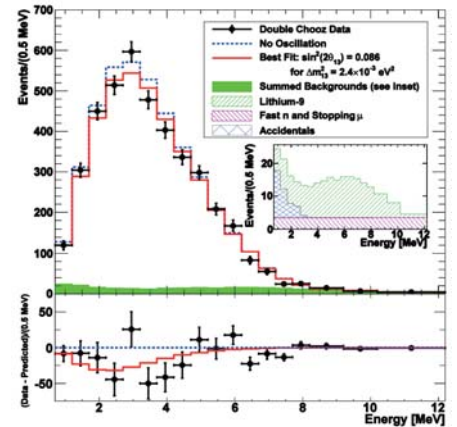


First Result

To extract the value of θ_{13} , a expected simulated spectrum was compared with the measured one, using the following χ^2 function:

$$\chi^2 = \left(N_i - \left(\sum_R N_i^{v,R} + \sum_b N_i^b(P_b) \right) \right) \times \left(M_{ij}^{\text{signal}} + M_{ij}^{\text{detector}} + M_{ij}^{\text{acc}} + \sum_b M_{ij}^b \right) \times \left(N_j - \left(\sum_R N_j^{v,R} + \sum_b N_j^b(P_b) \right) \right)^T + \sum_R \frac{N_R^{v,R}}{\sigma_R^2} + \sum_b \frac{N_b^b(P_b)}{\sigma_b^2}$$

which led to the result:



The final fit and first Double Chooz result is, therefore:

Rate Only fit: $\sin^2 2\theta_{13} = 0.104 \pm 0.030$ (stat.) ± 0.076 (sys.)

Rate & Shape fit: $\sin^2 2\theta_{13} = 0.086 \pm 0.041$ (stat.) ± 0.030 (sys.)

with no oscillation possibility excluded at 94.6%

High energy-resolution EELS and SXES studies on characteristic chemical shifts and charge transfer in Al-Si-Mn and Zn-Mg-Zr alloys

S. Koshiya (Physics, D3)

EELS: Electron energy-loss spectroscopy
SXES: Soft-X-ray emission spectroscopy

February 21 (Tue.), 2011.

In recent years, Hume-Rothery mechanism, which predicts a existence of a pseudogap around the Fermi level E_F , is accepted as a major reason for the stabilization of quasicrystals. The presences of pseudogap in quasicrystals were confirmed by X-ray photoemission spectroscopy and electron energy-loss spectroscopy (EELS). EELS experiments also pointed out characteristic chemical shifts in Al L-shell excitation spectra of Al-based quasicrystals, which suggested a decrease of valence electron charge at Al sites. Recently, a covalent bonding nature in quasicrystals was reported by MEM/Rietveld analysis. Thus, it is interesting to investigate the relation between a chemical shift and bonding nature of the quasicrystals.

Background1: Quasicrystals & approximant crystal

Quasicrystal
D. Shechtman et al., Phys. Rev. Lett. 53 (1984) 1951.

Approximant crystal
• Local structure is similar with a QC
• Characterized by degree of approximation
 $1/0 \rightarrow 1/1 \rightarrow 1/2 \rightarrow 2/3 \rightarrow \dots \rightarrow 1/n \ (n = (1+\sqrt{5})/2)$

QC's anomalous properties
• Lack translation symmetry
• High electronic resistance
• Hard and brittle

Quasicrystal and approximant crystal
→ Presence of pseudogap near the E_F
→ Complicated structure was stabilized??

Purpose of this study

Background2: EELS studies of Al-based quasicrystals

(a) $Al_{13}Cu_{15}V_{10}$

Al L-shell ($2p \rightarrow CB(3s, 3d)$)

Intensity (arb. units)

E_F

0.3 eV

0.2 eV

Energy loss (eV)

72 74 76

Cryst. Al-Cu-V
QC: Al-Cu-V
Am: Al-Cu-V

Diagram of core-level spectrum

Chemical shift

Fermi level

core level

V.B.

C.B.

Energy

Quasicrystal shows chemical shift
→ Decrease of valence electron charge
→ Presence of covalency??

Background3: MEM/Rietveld analysis

$1/0-Al_{12}Re$

$1/1-Al_{72.5}Re_{17.5}Si_{10}$

0.35 e \AA^3

0.35 e \AA^3

Covalent bonds

Interatomic bonds: Metallic

Interatomic bonds: Covalent

Covalent bonds between Al atoms
→ Decrease of valence electron charge??
→ Chemical shift (EELS)

The chemical shifts of $Al_{53}Si_{27}Mn_{20}$ alloys and Zn-Mg-Zr alloys were measured by using EELS and SXES.

- Investigate the relation between chemical shifts and crystalline order of $Al_{53}Si_{27}Mn_{20}$ alloys.
- Investigate the relation between chemical shift and bonding nature of the Zn-Mg-Zr (non-Al based) alloys.
- Those chemical shifts are compared with those of pure materials or their oxides.

High resolution EELS-TEM instrument

HR-EELS TEM
Acc. Volt: 60kV, 100kV
Analyzed area: ~1 μm^2
Energy resolution: 0.1 eV
Monochromator: Wien-filter
Analyzer: G-filter
Detector: IP (pixel size: 25 μm)

M. Mukai et al., Microsc. Microanal., 12 (Suppl. 2) (2006) 1206.

High resolution SXES-TEM instrument

SXES-TEM
Acc. Volt: 100kV, 200kV
Analyzed area: ~1 μm^2
Energy dispersion: 0.9 eV/pixel (Zn-L emission)
2.8 eV/pixel (Zr-L emission)
Grating: Multi-layer varied-line-space
Detector: CCD (pixel size: 13.5 μm)

M. Terauchi et al., J. Electron Microsc., 52 (2010) 251.

$Al_{53}Si_{27}Mn_{20}$ alloy (sample)

$Al_{53}Si_{27}Mn_{20}$ metal

Quench

Amorphous (Am)

5K/min. to 693K

Anneal

5K/min. to 803K

Quasicrystal (QC)

Crystal (Cryst)

A.P. Tsai et al., Phys. Rev. B, 49 (1994) 3569.
S. Koshiya et al., Philos. Mag., 91 (2011) 2309.

$Al_{53}Si_{27}Mn_{20}$ alloy (EELS & SXES spectra)

Electron diffraction pattern

BF image

Am

QC

Cryst

Al L-shell excitation: $2p \rightarrow CB(3s, 3d)$

Intensity (arb. units)

0.4 eV

0.1 eV

Energy loss (eV)

72 74 76

Cryst
Am
QC

(a) Al-K α : $2p \rightarrow 1s$ (1.5 keV)
(b) Si-K α : $2p \rightarrow 1s$ (1.7 keV)
(c) Mn-L: $3d \rightarrow 2p$ (640 eV)

10 eV

2 eV

0 eV

Intensity (arb. units)

QC
Am
Cryst

Zn-Mg-Zr alloy (composition & diffraction)

Primitive icosahedral quasicrystal (P-QC): $Zn_{83.5}Mg_{9.5}Zr_{7}$
Face-centered icosahedral quasicrystal (F-QC): $Zn_{75}Mg_{18}Zr_6$
1/1-approximant (1/1-AP) of F-QC quasicrystal: $Zn_{77}Mg_{19}Zr_4$

P-QC 2-fold axis

F-QC 2-fold axis

1/1-AP [100]

006

035

060

S. Ohnishi et al., Acta Mater., 57 (2009) 4727.

Zn-Mg-Zr alloy (EELS spectra)

Zn M-shell excitation: $3p \rightarrow CB(4s)$

Intensity (arb. units)

~0.2 eV

~0.4 eV

0.1 eV

Energy loss (eV)

84 85 86 87 88 89 90 91

1/1-AP
F-QC
P-QC
Zinc

Mg L-shell excitation: $2p \rightarrow CB(3s, 3d)$

Intensity (arb. units)

~0.14 eV

~0.19 eV

0.1 eV

Energy loss (eV)

46 47 48 49 50 51 52 53 54

1/1-AP
F-QC
P-QC
Pure Mg

Zn-Mg-Zr alloy (SXES spectra)

Zn L α : $3d \rightarrow 2p$ (1.0 keV)

Intensity (arb. units)

~2 eV

3.0 eV

Energy

1/1-AP
F-QC
P-QC
ZnO

Zr L α : $3d \rightarrow 2p$ (2.0 keV)

Intensity (arb. units)

~4 eV

10 eV

Energy

1/1-AP
F-QC
P-QC

• Chemical shifts are characteristic for not only Al-based quasicrystals but also Zn-Mg-Zr quasicrystals.

• All constituent atoms of quasicrystal show chemical shifts to larger binding energy side.
• The value of chemical shifts for quasicrystalline alloy comparable to their oxides.

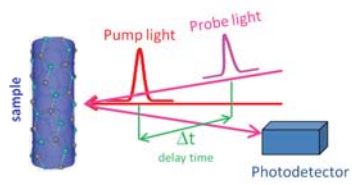
Summary

The chemical shift of $Al_{53}Si_{27}Mn_{20}$ and Zn-Mg-Zr alloys were investigated by using EELS and SXES.

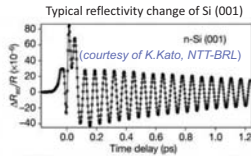
- Chemical shifts are characteristic for quasicrystalline materials of Al-based and Zn-Mg-Zr alloy.
→ Chemical shift is intrinsic for quasicrystalline state?
- All chemical shifts of quasicrystalline materials imply the decrease of valence electron charge.
→ Strongly suggest an increase of covalency in quasicrystalline states.
- The amount of shifts comparable to their oxides
→ Decreased valence electrons of quasicrystal are almost same to oxide?

Introduction to coherent phonon spectroscopy

Pump-probe spectroscopy

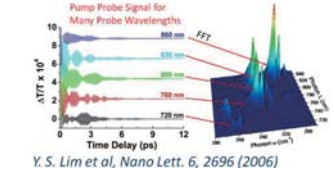
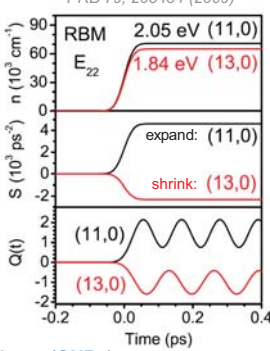
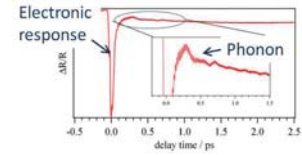


Either transmittance or reflectivity change can be measured



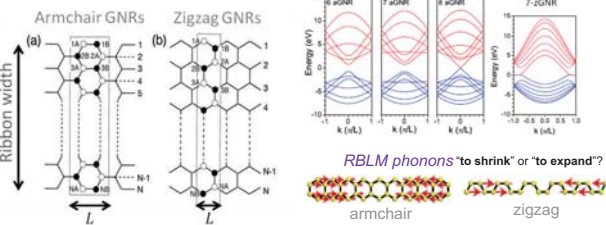
G. D. Sanders et al., PRB 79, 205434 (2009)

Coherent phonons in carbon nanotubes



Y. S. Lim et al. Nano Lett. 6, 2696 (2006)

Coherent phonons in graphene nanoribbons (GNRs)



Mechanism of coherent phonons

to find a general mechanism of coherent phonon generation in carbon nanotubes/nanoribbons

$$\frac{\partial^2 Q(t)}{\partial t^2} + \omega^2 Q(t) = S(t)$$

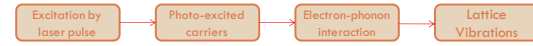
- What is an appropriate driving force for coherent phonons in carbon nanotubes?
- Why different nanotubes start coherent vibrations by uniquely expanding or shrinking their diameters?
- Role of electron-phonon and exciton-phonon interactions?

G. D. Sanders et al., PRB 79, 205434 (2009)

$$\frac{\partial^2 Q_{\beta}(t)}{\partial t^2} + \omega_{\beta}^2 Q_{\beta}(t) = S_{\beta}(t); \quad S_{\beta}(t) = -\frac{2\omega_{\beta}}{\hbar} \sum_{\mu\mu'} M_{\mu\mu'}^{\beta}(k) [f_{\mu\mu'}(k, t) - f_{\mu\mu'}^0(k)]$$

Q → CP amplitude Driving force el-ph interaction Carrier density

Generation of the coherent phonons:



Coherent phonon driving function

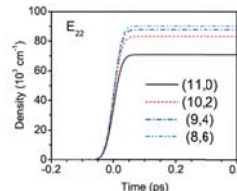
$$S(t) = \sum_{\mu\mu'} S_{\mu\mu'}(k) (f_{\mu\mu'}(k, t) - f_{\mu\mu'}^0(k))$$

Driving function kernel Photexcited carriers Equilibrium distribution

$$S_{\mu\mu'}^{\beta}(k) = -\frac{2\omega_{\beta}}{\hbar} [M_{\mu\mu'}^{\beta}(k) - M_{\mu\mu'}^{\beta}(k)]$$

Carrier density is obtained by solving Boltzmann equation

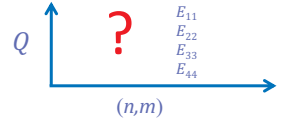
$$\frac{\partial f_{\mu\mu'}(k)}{\partial t} \Big|_{\text{gen}} = \frac{8\pi^2 e^2 u(t)}{\hbar n_0^2 \omega^2} \left(\frac{\hbar^2}{m_0} \right) \sum_{\mu''} [P_{\mu''\mu'}^{\mu}(k)]^2 [f_{\mu''\mu'}(k, t) - f_{\mu''\mu'}(k)] \delta(\Delta E_{\mu''\mu'}^{\mu}(k) - \hbar\omega)$$



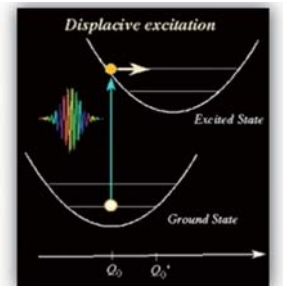
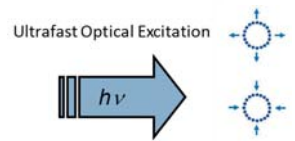
Purposes of the work

Map of initial lattice response

We can predict how a particular (n,m) nanotube/nanoribbons would respond to ultrafast excitation (shrink/expand)

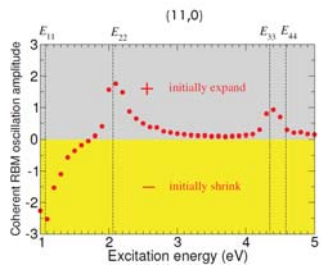


Calculation method



Results and Discussion

Excitation dependence (nanotubes)



"to expand" $M_{\text{el-ph}} < 0 \rightarrow S(t) > 0 \rightarrow Q(t) > 0$

"to shrink" $M_{\text{el-ph}} < 0 \rightarrow S(t) < 0 \rightarrow Q(t) < 0$

$$\frac{\partial^2 Q_{\beta}(t)}{\partial t^2} + \omega_{\beta}^2 Q_{\beta}(t) = S_{\beta}(t)$$

$$S_{\beta}(t) = \sum_{\mu\mu'} S_{\mu\mu'}^{\beta}(k) [f_{\mu\mu'}(k, t) - f_{\mu\mu'}^0(k)]$$

Magnitude of oscillations: $|Q| \propto |M_{\text{el-ph}}| |M_{\text{op}}|$

Eff. Mass theory: k-dependent el-ph interaction (nanotubes)

K. Sasaki et al., PRB 78, 235405 (2008)

ART Nugraha et al., PRB 84, 174302 (2011)

For zigzag tubes:

$$M_{\text{el-ph}} = \frac{A_0}{d_t} (-2g_{\text{off}} \cos \Theta(k))$$

For chiral tubes:

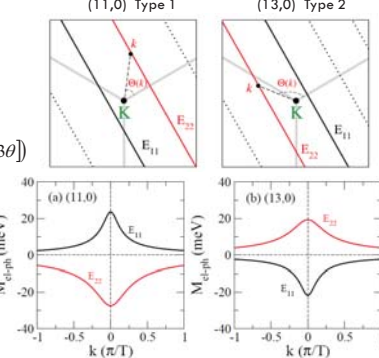
$$M_{\text{el-ph}} = \frac{A_0}{d_t} (-2g_{\text{off}} \cos[\Theta(k) + 3\theta])$$

to the right of the K-point

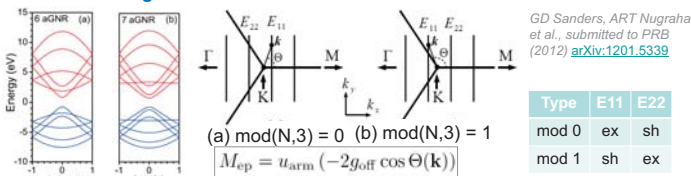
$$M_{\text{el-ph}} < 0 \text{ (expand)}$$

to the left of the K-point

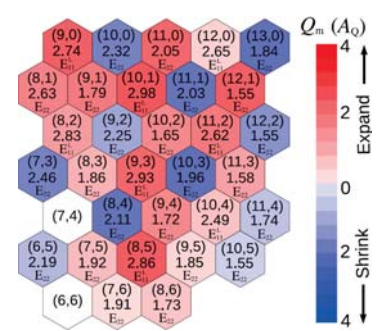
$$M_{\text{el-ph}} > 0 \text{ (shrink)}$$



Semiconducting armchair GNRs: mod 0 and mod 1



Chirality dependence and map of lattice response (nanotubes)

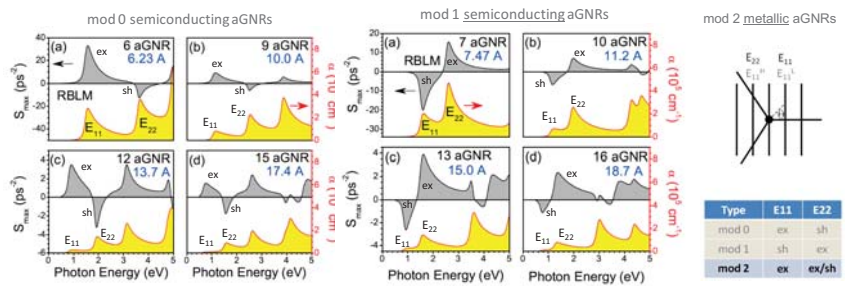


Deviated results from basic rules: (7,6), (9,5), (10,5) → near-armchair

Excitonic effects should also be considered in the near future.

Armchair nanotubes have a small coherent amplitude because of trigonal warping effect, where the density of states is split into lower and higher branches so that the two contributions cancel to each other.

Map of nanoribbon lattice response (armchair GNRs)



Summary and Conclusions

- Coherent RBM (nanotube) and RBM (nanoribbon) phonon starts by expanding or shrinking, depending on excitation energy and type/family
- Type dependence and excitation dependence of coherent phonon amplitudes originate from the k-dependent el-ph interactions

Laser Raman Spectroscopy and the *D* Band of Graphene

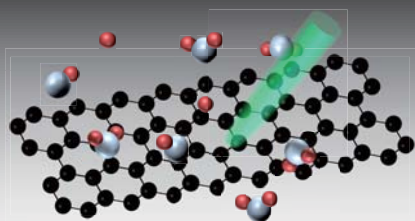
N. Mitoma¹, R. Nouchi² and K. Tanigaki^{1,2}

E-mail: mitoma@sspns.phys.tohoku.ac.jp

2-1-1 Katahira, Aoba-ku, Sendai 980-8577, Japan

Department of Physics, Tohoku University¹

WPI-Advanced Institute for Materials Research (WPI-AIMR), Tohoku University²



What is graphene?

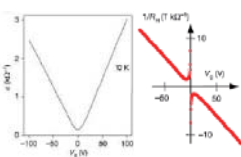


Optical microscope & AFM images

<http://www.gizmodo.com.au/2011/06/nokia-banks-on-wonder-material-graphene/>



Nobel prize in physics 2010

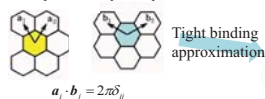


Carrier transport properties

K. S. Novoselov *et al.*, *Science* **306**, 666 (2004).
K. S. Novoselov *et al.*, *Nature* **438**, 197 (2005).

Energy band structure of graphene

Real space Reciprocal space



The expectation value of the energy

$$H\Psi = E\Psi \\ \Rightarrow \det(H - EI) = 0 \\ \therefore E = \hbar v_f |k|$$

$$H \approx \hbar v_f \begin{pmatrix} 0 & k_x - ik_y \\ k_x + ik_y & 0 \end{pmatrix}$$

Hamiltonian near K point



Effective mass is zero.

-High mobility ($\sim 200000 \text{ cm}^2 \text{ V}^{-1} \text{ s}^{-1}$) has been observed at LT.

K. I. Bolotin *et al.*, *Solid State Commun.* **146**, 351 (2008).

Background

Towards device application



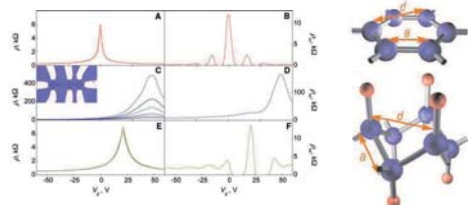
S. Bae *et al.*, *Nat. Nanotechnol.* **5**, 574 (2010).

Graphene's useful properties.

- Atomically thin ($\sim 0.34 \text{ nm}$), i.e., transparent and flexible material.
- High charge-carrier mobility due to its characteristic energy band.
- Eco-friendly material, because all of it is constructed from carbon atoms.

Problems for device application.

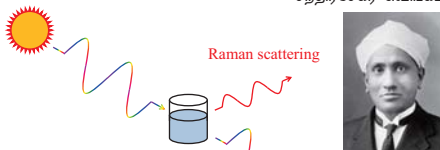
- Graphene cannot be applied as a switching device due to its no band gap...



D. C. Elias *et al.*, *Science* **323**, 610 (2009).

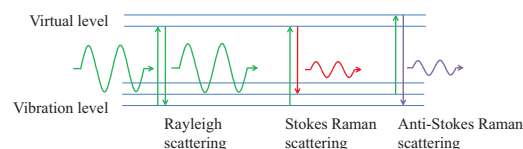
Raman spectroscopy

சந்திரசேகர வெங்கடராமன்



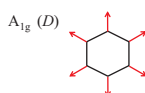
Rayleigh scattering

C. V. Raman and K. S. Krishnan, *Nature* **121**, 501 (1928).

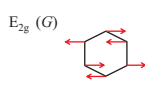


Water molecules are found to exhibit nonnegligible reactivity with graphene under strong light irradiation. The reaction is triggered by laser Raman spectroscopy measurement itself, and the *D* band (*ca.* 1340 cm^{-1}) becomes larger as the laser irradiation is prolonged. The electronic transport properties of the graphene derivative are also measured and the mobility is found to be reduced. These results could originate from modification of graphene with hydrogen.

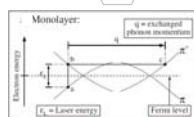
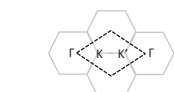
Raman reactive modes



This vibration mode is forbidden on a "perfect" graphene lattice.



Can be seen on all Graphene-related materials.



Double resonant Raman process, the origin of the 2D band.

A. C. Ferrari *et al.*, *Phys. Rev. Lett.* **97**, 187401 (2006).

Time evolution of the *D* bands from 3 types of samples.

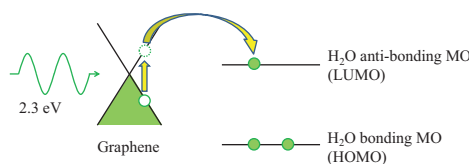
$$I_{\text{Raman}}(D) \propto n_{\text{defects}}$$

$$n_{\text{defects}} = N \{ 1 - \exp[-\sigma_{\text{H}_2\text{O}} \Phi] \}$$

Reaction cross section Incident photon flux

- Graphene on a dry substrate exhibited a low reactivity.
- The reaction probability is expressed by the number of water molecules, reaction cross section, and incident photon flux. These are the experimental parameters.
- The reaction probability of graphene on a wet substrate converged to similar value compared to that of on a dry substrate.
- The negative slope appeared on a Au-deposited graphene sample implies that the material went into the amorphousization trajectory.

Why the reaction is triggered by laser irradiation?



Decomposition of water molecules is accelerated by the photoradiation.

Charge-carrier transport property

- Charge-carrier mobility reduced after the light irradiation. This is thought to be the result of hydrogenation of the graphene plane.
- The resistivity at Dirac point slightly increased after the irradiation.

Summary & Future tasks

- We observed the enlargement of the *D* band as the laser irradiation was prolonged.
- The peak position shifts of the *D*, *G*, and *2D* are also observed. These indicate the carrier doping to the graphene plane.
- A study on electronic transport properties is still lacking. We need to measure at a variety of temperatures to clarify the carrier scattering mechanism.
- A study on the photochemical reactivity of multi-layer graphene is needed.

Keywords:

Graphene,

Raman spectroscopy,

Photochemical reaction,

Charge-carrier mobility

Spectral-Function Sum Rules in Supersymmetry Breaking Models

Mistutoshi Nakamura

Collaborated with Ryuichiro Kitano, Masafumi Kurachi, and Naoto Yokoi
arXiv:1111.5712 [hep-ph]

Abstract:

If supersymmetry is realized in nature, it must be spontaneously broken. We derive relations among physical quantities as a consequence of supersymmetry, which generally applies in models of dynamical supersymmetry breaking.

Symmetry

Symmetry is a field transformation which does not change action.

$$\phi \rightarrow \phi + \delta\phi$$

$$S[\phi] \rightarrow S[\phi + \delta\phi] = S[\phi]$$

Supersymmetry

Supersymmetry is symmetry that transforms from boson to fermion and vice versa.

$$\delta_Q(\text{boson}) = \text{fermion}$$

$$\delta_Q(\text{fermion}) = \text{boson}$$

Spontaneous Symmetry Breaking

$$\langle \delta(\mathcal{O}_1 \dots) \rangle = 0$$

corresponds to unbroken symmetry

→ Conservation of charge

$$\langle \delta(\mathcal{O}_1 \dots) \rangle \neq 0$$

corresponds to broken symmetry

→ Massless Nambu-Goldstone particle

Sum rules

S. Weinberg, (1967).

$$\frac{\langle \mathcal{O} \rangle}{q^n} \sim \langle \delta(\mathcal{O}_1 \mathcal{O}_2) \rangle \neq 0 \quad \text{Spectral function}$$

$$= - \int_0^\infty d\sigma^2 \frac{\rho(\sigma^2)}{q^2 - \sigma^2 + i\epsilon}$$

$$= \frac{c_1}{q^2} + \dots + \frac{c_{n/2}}{q^n} + \dots$$

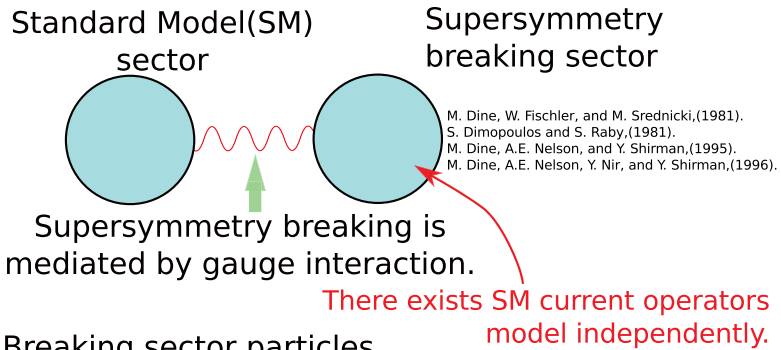
$$\underbrace{\frac{c_1}{q^2}}_{=0} \dots \quad c_n = \int_0^\infty d\sigma^2 \rho(\sigma^2) \sigma^{2(n-1)}$$

→ One can derive sum rules which do not depend on details of underlying model.

This n is determined by the dimension of \mathcal{O} . We denote this dimension as d .

→ We apply sum rule to supersymmetry breaking model.

Gauge Mediation



Breaking sector particles

	mass	decay constant
Spin-0	m_0	f_0
Spin-1/2	$m_{1/2}$	$f_{1/2}$
Spin-1	m_1	f_1

By taking $\mathcal{O}_1, \mathcal{O}_2$ as currents of SM gauge group, we obtain

$$f_0 = f_{1/2} = f_1 \equiv f_h.$$

By using this scalar and gaugino mass are computed to be

$$m_s^2 = \frac{g^4 c_2 f_h^2}{(4\pi)^2} \log \frac{m_0^2 m_1^6}{m_{1/2}^8} \quad m_\lambda = \frac{g^2 f_h^2}{m_{1/2}}$$

→ More generally there exists supercurrent operator.

Supercurrent

Supercurrent is a current corresponding to supersymmetry.

→ We take supercurrent as $\mathcal{O}_1, \mathcal{O}_2$.

	mass	decay constant
Scalar		$f_\phi \quad c_\phi$
Pseudoscalar		$f_\pi \quad c_\pi$
Massless spin-1/2	0	$f \quad f'$
Massive spin-1/2	m_χ	f_χ
Massive spin-1	m_v	f_v
Massive spin-3/2	m_ψ	f_ψ
Massive spin-2	m_h	m_P

Sum rule for $d_0 = 3$ and 4

$$|f'|^2 + |f_\chi|^2 + \frac{2}{3}|f_\psi|^2 = f_\phi^2 + \frac{8}{3}m_P^2,$$

$$f_\phi^2 + \frac{8}{3}m_P^2 = f_\pi^2 + f_v^2,$$

Sum rule for $d_2 = 4$

$$f^2 f' = m_\psi f_\psi^2 = -\frac{3}{4}m_\chi f_\chi^2,$$

$$f_\phi c_\phi^2 = 0, \quad f_\pi c_\pi^2 = 0,$$

Terahertz pulse shaping via difference frequency mixing of shaped optical pulses

Koji Uematsu^{a,b} and Chiko Otani^{a,b}

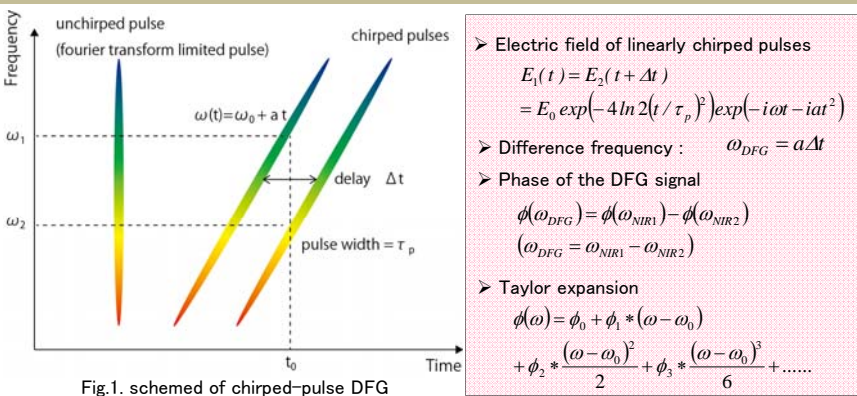
Department of Physics, Tohoku University ^a Terahertz Sensing and Imaging Laboratory, RIKEN ^b

The 4th GCOE International Symposium on "Weaving Science Web beyond Particle-Matter Hierarchy"

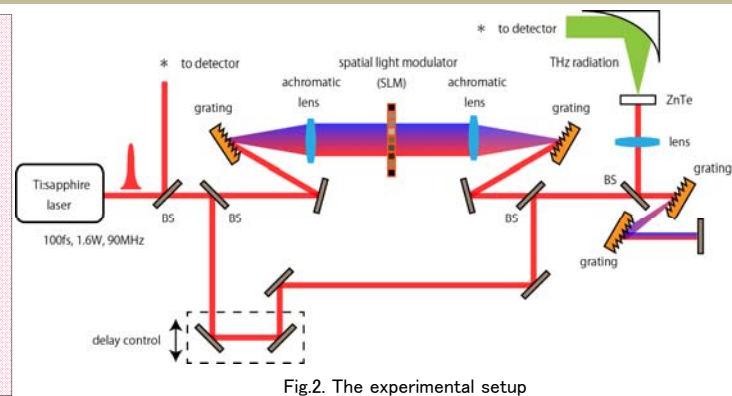
Abstract

Femtosecond pulse shaping techniques which exist in wavelength range from visible to mid-infrared have been developed for various applications, such as coherent control of matter and chemical reactions. In terahertz (THz) region, the use of intense sub-picosecond THz pulses provides tools in non-linear and time-domain THz spectroscopy, but no efficient way for THz pulse shaping exists. Complex and precise control techniques of THz temporal waveforms are desirable for those applications. In this poster, we propose and demonstrate a novel method of arbitrary pulse shaping. It is based on spectral phase transfer to THz pulses by a difference-frequency generation technique.

Principle – difference-frequency generation (DFG)



Experimental setup



Experiments and Results

the phase measurement of a shaped ultrashort optical pulse : spatial spectral interferometry (SSI)

Fig.3. SSI setup

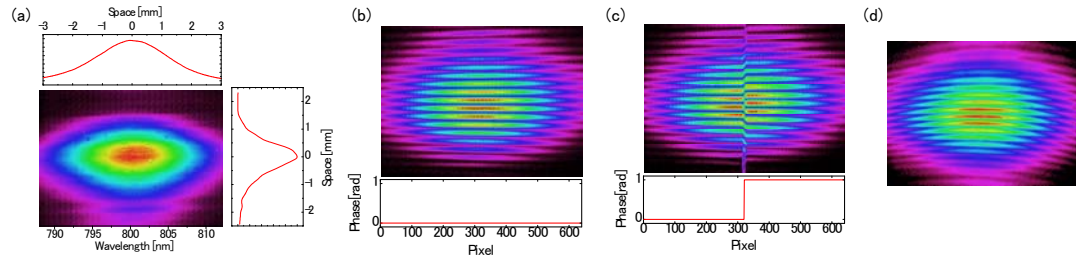
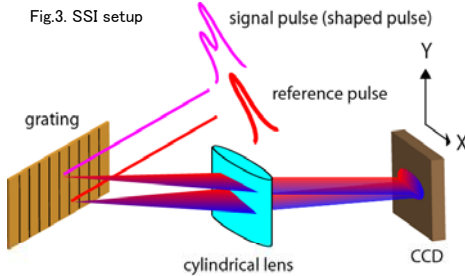


Figure 4. (a) Beam profile of the reference pulse measured by a CCD and SSI fringe pattern with (b) no modulation (c) phase change of π at pixel 320 (d) second-order chirp.

Signal : $I_{sig}(r, t) = E_{sig}(r) \exp[i(k_{sig} \cdot r + \phi_{sig}(r) - \omega t)]$
 Reference : $I_{ref}(r, t) = E_{ref}(r) \exp[i(k_{ref} \cdot r + \phi_{ref}(r) - \omega t)]$
 The fringe pattern on the CCD plane
 $I(x, y) = |I_{ref}(x, y, t) + I_{sig}(x, y, t)|^2 = E_{ref}^2(x, y) + E_{sig}^2(x, y) + 2E_{ref}(x, y)E_{sig}(x, y) \cos[2k_y \sin \theta + (\phi_{ref}(x, y) - \phi_{sig}(x, y))]$

z waveforms for chirped-pulse

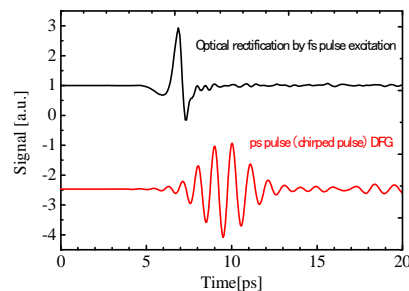


Figure 5. THz waveforms for femtosecond pulse excitation (black) and chirped pulse DFG (red).

- Femtosecond pulse width 100fs
- Group delay dispersion of chirped pulse -0.2ps^2

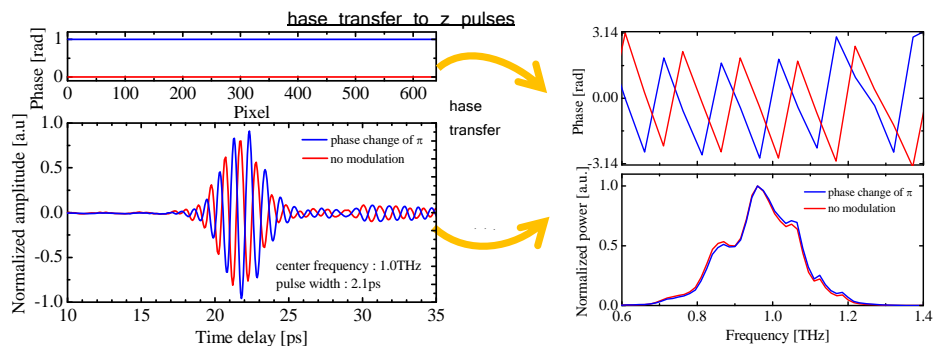


Figure 6. (left) THz waveforms for no modulation and phase change of π at S M (right) phase and frequency spectra of the radiation.

Spectral phase control at S M

$$\phi(\omega) = \phi_0 + \phi_1 * (\omega - \omega_0) + \phi_2 * \frac{(\omega - \omega_0)^2}{2} + \phi_3 * \frac{(\omega - \omega_0)^3}{6} + \dots$$

z_pulses

Tailored THz waveforms for various applications

Summary

- I have demonstrated that the phase profile of a shaped pulse in the near-infrared is transferred to a pulse in the terahertz via a difference frequency mixing (DFM) process.
- This phenomenon will allow THz waveform shaping by manipulating the high-order dispersion of chirped pulse.
- The THz pulse shaper can be applied to the study of intermolecular vibrations, structure of large molecules and molecular rotational level and the control of the coherent phonons, antiferromagnetic spin waves and chemical reactions.

On the development of resonant inelastic x-ray scattering for high-pressure experiments

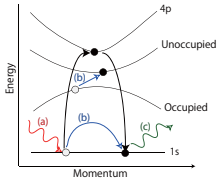
M. Yoshida

Dept. of Phys. Tohoku Univ. / JAEA SPring-8

Background

Pressure is one of the most important external parameters along with temperature and magnetic field. By inducing pressure to a material, we can control its physical properties through shrinkage of the crystal lattice. While investigating the electronic structure is central to the understanding of physical properties, such studies are still hardly carried out due to the stringent conditions imposed by the use of high-pressure cells.

Why RIXS?



<advantages>

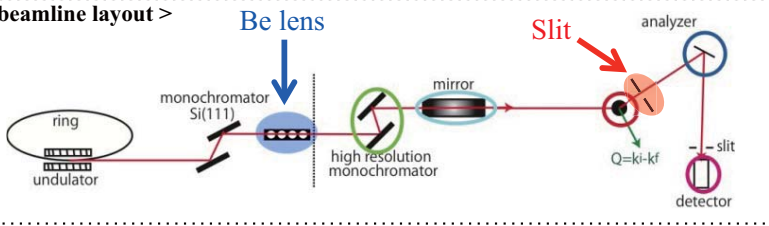
- momentum dependence
- bulk sensitivity
- element sensitivity
- **photon-in-photon-out spectroscopy**

Applied to high-pressure experiments

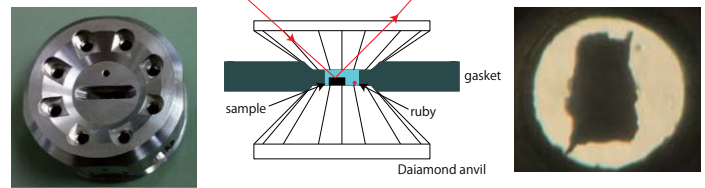
Experiments

✓ We performed RIXS experiments at BL11XU, SPring-8

< beamline layout >



DAC



- diamond-in-diamond-out configuration
- diamond size
h1 = 1 mm, h2 = 1.5 mm
(asymmetry thickness)
C1 = 1 mm (for large sample volume)
- sample size 400×250×70 μm³
- fluorinert
- SUS301
- pressure determination for ruby fluorescence

Intensity is reduced due to diamond, about two or three hundredth parts.

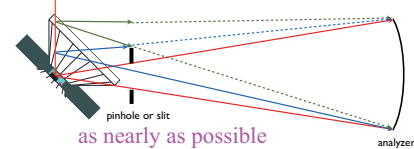
✓ We introduced two optical components for high-pressure RIXS.

- Be lens

For focusing x-rays along vertical direction

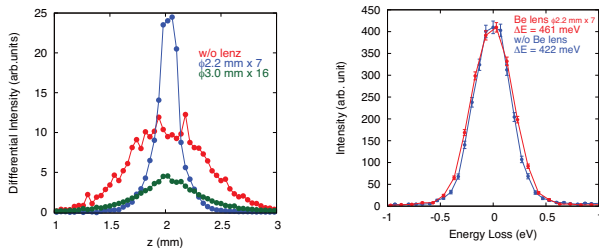
- Slit located after sample

For removing the scattering x-ray from diamond



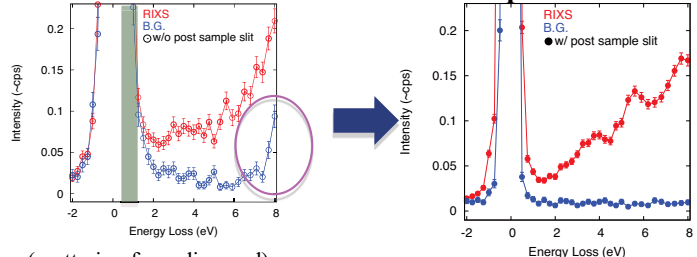
Estimation of two optical components

< Be lens >



-Using $\phi 2.2 \text{ mm} \times 7$, vertical beam size reduce by half but peak height become twice and energy resolution does not change so much.

< Slit located after sample >

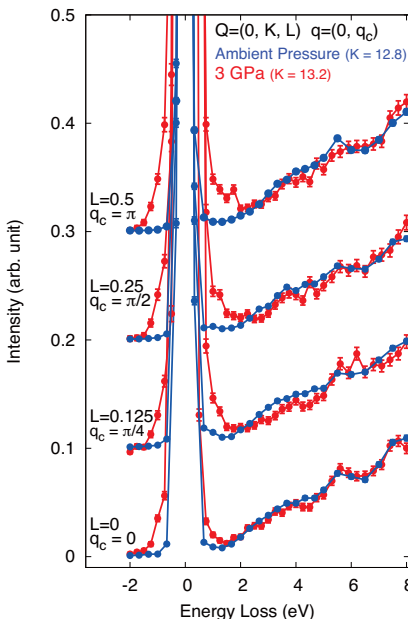


(scattering from diamond)

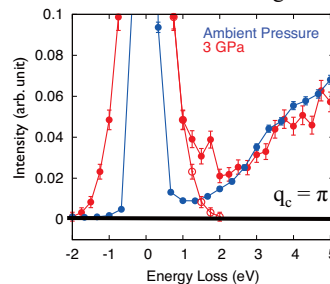
- asymmetry of elastic scattering
- some excitation around 8 eV

Suppress scattering from diamond !!

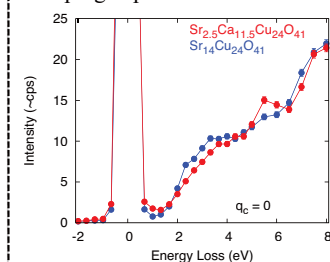
Results



focus on intraband region



cf. doping dependence



✓ In order to compare the data measured at ambient pressure, we normalize them using the integrated intensity between 5 – 8 eV. (scale factor : ~ 1/160)

We have succeeded to obtain Q-resolved RIXS spectra under high pressure for the first time though intensity become very weak.

- crossover around 2 eV energy loss
→ Similar behavior have shown in doping dependence.
- large enhancement of the intensity in the intraband region
→ It is considered that a number of holes in the ladders increases.

Summary

- ✓ We succeeded to develop RIXS method for high-pressure experiments due to the introduction of Be lens and a slit located after sample.
- ✓ We obtained Q-resolved RIXS spectra under high pressure for the first time and observed the change of them in pressure.
- ✓ In order to understand the effect of physical pressure in detail, we have to measure RIXS spectra under several pressure.

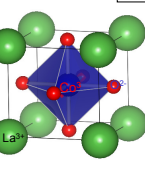
Acknowledgements

This work was supported by the follows. I would like to thank them.

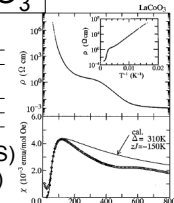
K. Ishii, I. Jarrige, T. Watanuki, K. Ohwada, K. Tsutsui, J. Mizuki, (JAEA, SPring-8), Y. Murakami (IMSS, KEK), N. Hiraoka, H. Ishii, K. D. Tsuei (NSRRC), A. Q Baron (RIKEN), T. Tohyama (YITP, Kyoto Univ.), S. Maekawa (ASRC, JAEA), K. Kudo (Okayam Univ.), Y. Koike (Tohoku Univ.), K. Kumagai (Hokkaido Univ.), Y. Endoh (IIAS)

1. Introduction

Perovskite-type cobalt oxide LaCoO₃

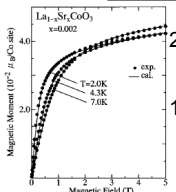


Co³⁺ (3d⁶)
e_g
t_{2g}
Low-spin (LS) state (S = 0)
Intermediate-spin (IS) state (S = 1)
High-spin (HS) state (S = 2)

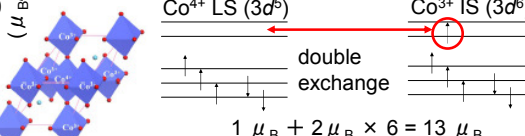


- Pseudo-cubic crystal structure.
- The Co³⁺ ion exhibits spin-state degree of freedom: LS state, IS state, and HS state.
- Below 100 K, non-magnetic insulator state appears with LS state [1].

Hole doped system La_{1-x}Sr_xCoO₃



• Huge magnetic moment of over 10 μ_B /hole is observed below 10 K [1].



- To explain the huge magnetic moment, spin-state polaron model comprising seven Co ions was proposed [2, 3].
- The origin of spin polaron is considered to be double exchange interaction between e_g orbitals [4].

Electron doped LaCoO₃

1. La_{1-x}Ce_xCoO₃ (Ce⁴⁺)

- Above x > 0.05, single-phase sample is difficult to synthesize [5].
- Few data have been reported.

2. La_{1-x}Te_xCoO₃ recently reported (Te⁴⁺ [6] or Te⁶⁺)

- It is easy to synthesize up to x = 0.25 [6].

Which is correct, Te⁴⁺ or Te⁶⁺ ?

2. Purposes

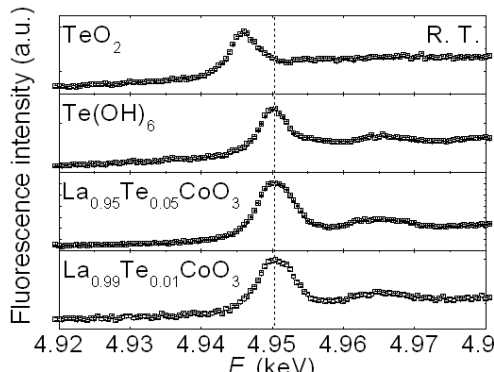
- Determination of the valence number of doped Te ion.
- Experimental examination of the spin-state polaron in electron-doped system.

3. Experiments

- Polycrystalline samples La_{1-x}Te_xCoO₃ (x = 0.01, 0.05) were synthesized. (1050 °C for 24 h and 1050°C for 75 h)
- X-ray fluorescence (XRF) analysis method at PF BL-4C, Tsukuba, Japan. The fluorescence energy range from 3.6 to 4.2 keV was selected in energy dispersive XRF spectrum. We used TeO₂ and Te(OH)₆ as reference compounds of Te⁴⁺ and Te⁶⁺, respectively.
- Magnetization measurements down to 1.8 K in magnetic field up to 7 T at Center for Low Temperature Science, IMR, Sendai, Japan.

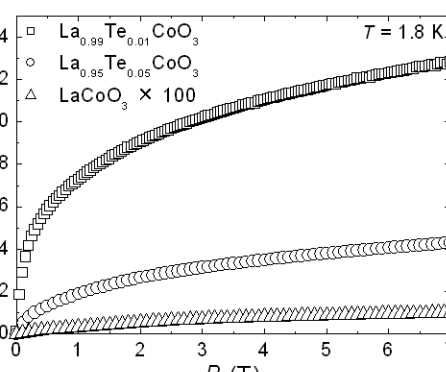
4. Results

Fluorescence spectra around Te-L₁ edge



- It is considered that the valence number of doped Te ion is six.

Magnetization of La_{1-x}Te_xCoO₃



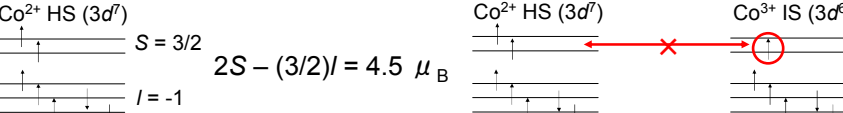
- Magnetization of La_{0.99}Te_{0.01}CoO₃ reaches 13 μ_B /Te at 7 T.

5. Discussion

- The valence number of doped Te ion: Te⁶⁺

Each Te ion introduces three electrons.

13 μ_B /Te \rightarrow 4.3 μ_B /electron consistent with the Kanamori theory [7].



6. Summary

- The valence number of doped Te ion: Te⁶⁺
- It is considered that the spin-state polaron does not realize in the lightly Te-doped LaCoO₃.

7. Future

- Growth of single crystalline La_{1-x}Te_xCoO₃
- Construction of a Te concentration x – T phase diagram
- Investigation of thermoelectric properties

References

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[2] D. Phelan *et al.*, Phys. Rev. Lett. **97**, 235501 (2006).

[3] A. Podlesnyak *et al.*, Phys. Rev. Lett. **101**, 247603 (2008).

[4] D. Louca *et al.*, Phys. Rev. Lett. **91**, 155501 (2003).

[5] J. Kirchnerova *et al.*, Appl. Catal. A:Gen **231**, 65 (2002).

[6] G. H. Zheng *et al.*, J. Appl. Phys. **103**, 013906 (2008).

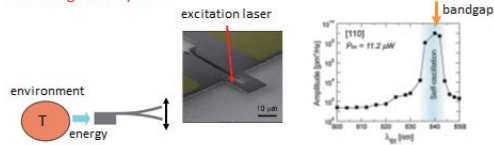
[7] J. Kanamori, Prog. Theor. Phys. **17**, 177 (1957).

Laser cooling of GaAs/AlGaAs cantilever by exciton related GaAs optical absorption

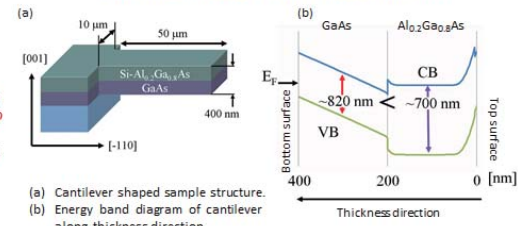
Takayuki Watanabe

Background

Vibration control with carrier induced optomechanical coupling
 → Large amplification occur at bandgap absorption
 → This technique is sensitive to the spectrum of bandedge structure
 ↓
 In order to characterize fine absorption spectrum with our technique, we used a GaAs/AlGaAs heterostructured cantilever which have sharp bandedge absorption.

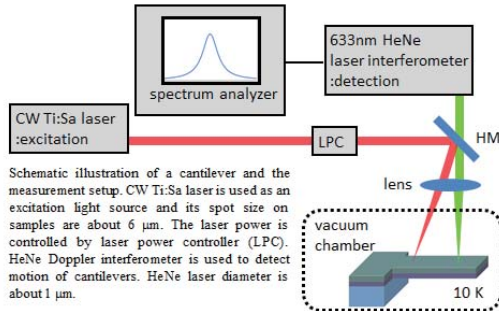
H. Okamoto et al, Phys. Rev. B **84** (2011) 014305

Sample structure and mechanism



(a) Cantilever shaped sample structure.
 (b) Energy band diagram of cantilever along thickness direction.
 Resonance frequency of a cantilever is ~114 kHz. Quality factor is ~1300 at 10 K.

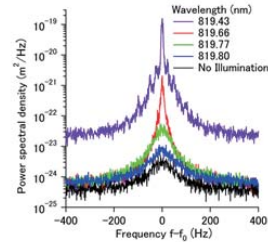
Measurement setup



Schematic illustration of a cantilever and the measurement setup. CW Ti:Sa laser is used as an excitation light source and its spot size on samples are about 6 μm. The laser power is controlled by laser power controller (LPC). HeNe Doppler interferometer is used to detect motion of cantilevers. HeNe laser diameter is about 1 μm.

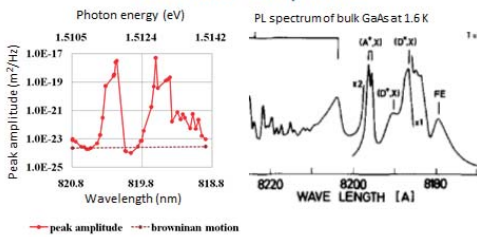
Experiment 1 optical absorption measurement of GaAs around exciton state with micro-scale mechanical cantilever

Cantilever's PSD with excitation laser



Excitation laser's wavelength dependence of cantilever's power spectral density (PSD) shows strong interaction when the wavelength is just below of the GaAs bandgap at 10 K. The resonance peak becomes sharper and larger when the excitation laser wavelength is slightly reduced.

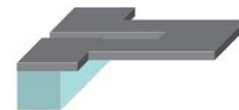
Wavelength dependence of cantilever's amplitude



Peaks which appear in the wavelength dependence of cantilever's peak amplitude shows good agreement to the typical PL peak of bulk GaAs.

Experiment 2 mechanical mode cooling of cantilever with cavity-free opt-electro-mechanical coupling

Sample structure

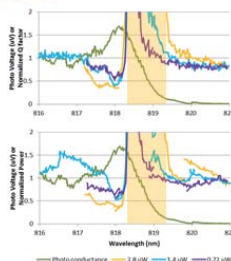


Cantilever structure (sample2)
 20 μm long, 14 μm width, 400 nm thickness (including 4 μm wide and 10 μm long legs) and [110]-oriented
 Resonance property at 10 K
 $R_F = 383.1 \text{ kHz}$
 $Q \sim 10000$

Measurement setup is the same of experiment1.

Laser power dependence of mechanical response

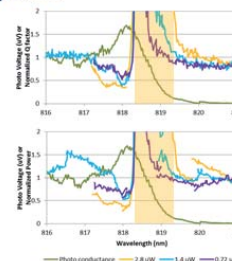
Wavelength dependence of flat noise added mechanical vibration was measured in detail. The added flat noise level was 18 μVrms/rHz.



Cooling like effect was observed around 818 nm.

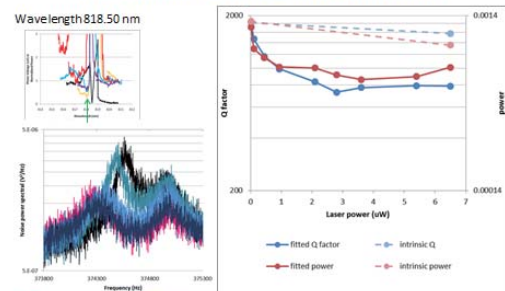
Laser power dependence of mechanical response

Wavelength dependence of flat noise added mechanical vibration was measured in detail. The added flat noise level was 18 μVrms/rHz.



Cooling like effect was observed around 818 nm.

Mechanical resonance peak de-amplification (flat noise is added for actuation)



Summary

- We measured the excitation laser wavelength dependence of thermo-mechanical vibration of GaAs/AlGaAs cantilever.
- A sensitive wavelength dependence of thermal vibration was observed.
- A significant vibration amplification was confirmed around exciton-related PL peaks.
- A mechanical mode cooling is observed around GaAs exciton peak.

Role of non-collective excitations in heavy-ion reaction around the Coulomb barrier

Shusaku Yusa

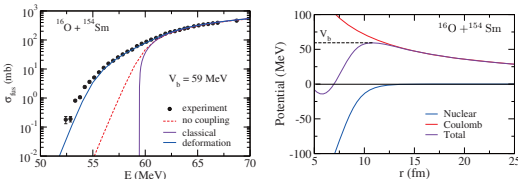
Department of Physics, Tohoku University

Introduction heavy-ion reaction

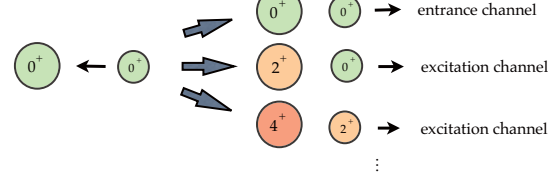
nucleus is composed of protons and neutrons
 ⇒ excitation during the collision

important in reactions near the Coulomb barrier
 quantum tunneling amplify the effect

e.g. large enhancement of fusion cross sections at sub-barrier energies



Coupled-channels method



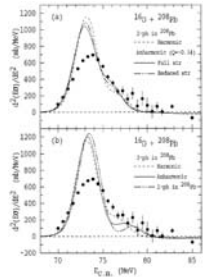
Quantum theory which takes into account nuclear excitations

conventionally, a few **low-lying collective excitations** (rotation, vibration) are taken into account

successfully accounted for heavy-ion fusion reaction and quasi-elastic scattering (sum of elastic, inelastic and transfer reaction) experiments

¹⁶O + ²⁰⁸Pb

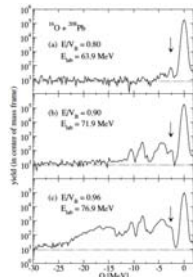
fusion reaction (barrier distribution)



C. R. Morton *et al.*, PRC60, 044608(1999)

failed to account for the data by the C.C. analysis with only collective excitations

quasi-elastic scattering $\theta_{lab} = 163^\circ$ (Q-value distribution)



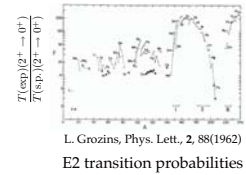
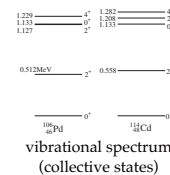
M. Ever *et al.*, PRC78(2008)034614

increasing contribution from higher energy region (non-collective excitations)

Can non-collective excitations improve the agreement?

Nuclear excited states

	collective states	non-collective states
number of nucleons	many nucleons	a few nucleons
appearance of energy levels	low-lying, regular and high-lying (giant resonances)	relatively high excitation energy, irregular
transition strength	dozens of times larger than single-particle transition	comparable to single-particle transition



Coupled-Channels method

coupled-channels equations (in the isocentrifugal approximation):

$$\left[\frac{d^2}{dr^2} + k^2 - \frac{J(J+1)}{r^2} - \frac{2\mu}{\hbar^2} (V_C(r) + V_N(r) + \epsilon_n) \right] u_n^J(r) = \sum_m \frac{2\mu}{\hbar^2} V_{nm}(r) u_m^J(r)$$

$k = \sqrt{\frac{2\mu E}{\hbar^2}}$ ϵ_n : excitation energy for the n -th channel

$V_{nm}(r)$: coupling matrix element (induce excitations)

collective excitations

²⁰⁸Pb: 3: 2.615 MeV, $\beta_3 = 0.122$

5: 3.198 MeV, $\beta_5 = 0.058$

2⁺: 4.085 MeV, $\beta_2 = 0.058$

¹⁶O: 3: 6.13 MeV, $\beta_3 = 0.733$

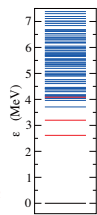
non-collective excitations

70 states in ²⁰⁸Pb up to 7.382 MeV

deformation parameter β_3

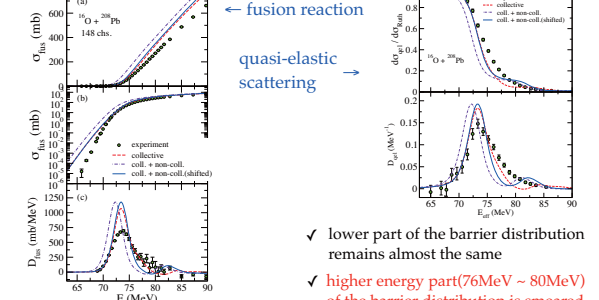
excitation energy ϵ^*

²⁰⁸Pb: high precision proton inelastic scattering experiment (W.T. Wagner, *et al.*, PRC12 757(1975))



Results ¹⁶O + ²⁰⁸Pb effect of non-coll. excitations

$$D_{qel} = -\frac{d}{dE} \left(\frac{\sigma_{qel}(E, \pi)}{\sigma_R(E, \pi)} \right)$$



✓ fusion cross sections are not much altered (only shift in energy)

✓ lower part of the barrier distribution remains almost the same

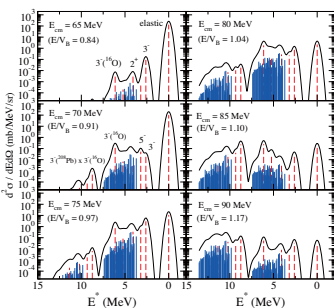
✓ higher energy part (76 MeV ~ 80 MeV) of the barrier distribution is smeared

✓ agreement with the data is not improved (large difference in main peak)

¹⁶O + ²⁰⁸Pb Q-value distribution for quasi-elastic scattering

Energy dependence of Q-value distribution ($\theta = 170^\circ$)

red - collective channels, blue - non-coll. channels



✓ for lower incident energy, collective channels are dominant

✓ as the incident energy increases, non-collective excitations become important

✓ energy dependence is consistent with the experimental Q-value distribution

Summary

heavy-ion reaction with non-collective excitations
 → application to ¹⁶O + ²⁰⁸Pb system

- ✓ information for non-collective excitations of ²⁰⁸Pb has been obtained from high resolution proton-inelastic scattering experiment
- ✓ higher part of the barrier distribution is smeared due to the non-collective excitations (both fusion and quasi-elastic scattering)
- ✓ agreement with the experiment is not improved
- ✓ energy dependence of the Q-value distribution is consistent with the experiment

Future perspectives

- application to ²⁰Ne + ^{90,92}Zr systems (importance of non-coll. excitations is suggested)
- description non-collective excitations with random matrix theory (for the purpose of application to the system whose information for non-coll. states is absent)

No.31 Improvement of the detector systems for Neutral Kaon Spectrometer 2

Physics D1 Fumiya Yamamoto

✓ Introduction

- Measurement of meson photoproduction with Neutral Kaon Spectrometer 2 (NKS2)
 - Elementary photoproduction of strangeness : $\gamma n \rightarrow K^0 \Lambda$
 - Non-quasi-free photoproduction on the deuteron : $\gamma d \rightarrow \pi^+ \pi^- p n$ ($\gamma d \rightarrow \Delta^{++} \Delta^-$)
 - a consistent cross section at wide photon energy range 0.54 - 1.25 GeV (Fig.1).

• Improvement of NKS2

- Lead glass Cherenkov counter (EVLG, Fig2).
 - Decrease trigger rate of e^+e^- background \rightarrow High beam intensity.
- Vertex Drift Chamber (VDC, Fig3).
 - Vertex resolution \rightarrow Higher the momentum and angular resolution.
 - Acceptance (Fig.4) : K^0 3 times, Λ 1.5 times, $K^0 \Lambda$ 8 times.
 - Elimination of a Fermi motion correction of neutron in deuteron by coincidence $K^0 \Lambda$.
- Development of new tagger with MPPC.
 - Photon energy range : 0.54 - 1.1 GeV \rightarrow 0.54 - 1.25 GeV.
 - Photon energy resolution : 6 MeV \rightarrow 3 MeV.

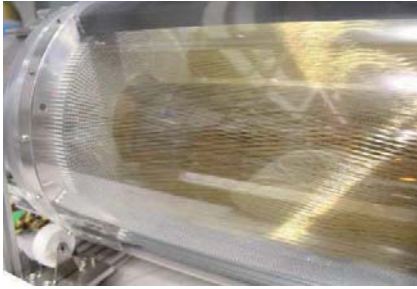


Fig.3 : Photograph of VDC. This counter located close to a target. The number of channel is 626. Three dimensional analyses are possible by wire having tilt for z-axis.

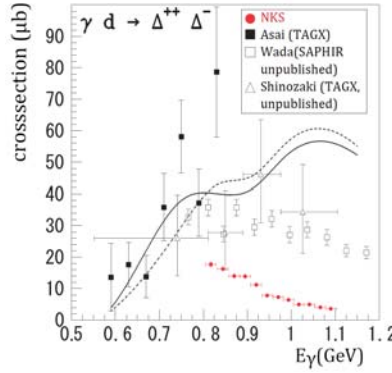


Fig.1: The Cross section for the $\gamma d \rightarrow \Delta^{++} \Delta^-$.
Ref. K. Hirose, et al., Phys. Lett. B 674 (2009) 17.

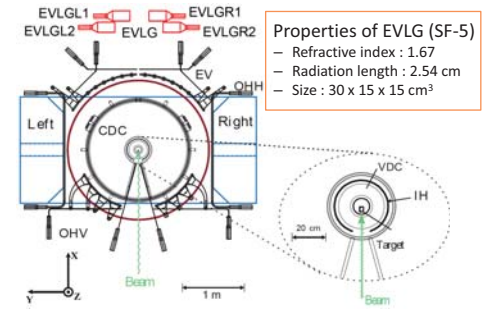


Fig.2: NKS2 scheme and position and properties of EVLG. NKS2 consists of dipole magnet (a diameter 80 cm, a gap 68 cm), hodoscopes, drift chamber and liquid deuterium target on center of magnet. EVLG located on beam plane, 3 m from center of magnet.

• Report in this poster

- Evaluation of a quality of trigger level e^+e^- background reduction with ADC of EVLG.
- Evaluation of DAQ efficiency before and after incorporating EVLG as veto in the trigger.

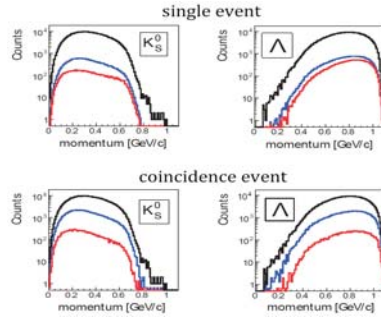


Fig.4 : Compare of a single event K^0 and Λ and a coincidence event $K^0 \Lambda$ with / without VDC by Monte-Carlo simulation. Black line : Total event. Red line : Acceptance of NKS2 without VDC. Blue line : Acceptance of NKS2 with VDC.

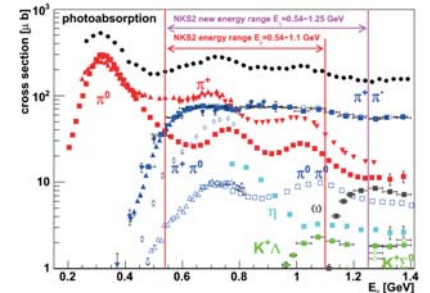


Fig.5 : The cross section for photoreaction on proton target. Red line : previous tagged photon energy range 0.54-1.1 GeV. Purple line : new tagged photon energy range 0.54-1.25 GeV.

✓ Particle identification

- Particle selection by its mass :

$$m^2 = p^2(\beta^2 - 1)$$

m : mass
 p : momentum
 β : velocity
- Condition of π selection :

$$-0.5 < m^2 < 0.5$$

$$|p| > (0.144/(\beta^{-1}-0.2)-0.8)$$

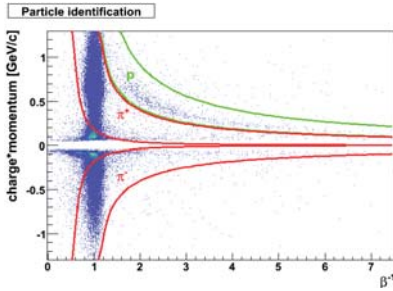


Fig.6 : Particle identification, β inverse vs momentum. Green line : proton. Red line : π^+ and π^- .

✓ Results

- Electron reduction : $\{N_{\text{threshold}}(\text{without}) - N_{\text{threshold}}(\text{with})\} / \{N_{\text{all}}(\text{without}) - N_{\text{all}}(\text{with})\}$
 π loss : $N_{\text{threshold}}(\text{with}) / N_{\text{all}}(\text{with})$
- Effect of EVLG veto.
 DAQ efficiency : Improvement about 10% at beam rate 1.5 MHz.
 Beam intensity : Improvement about 30% at DAQ efficiency 0.8.

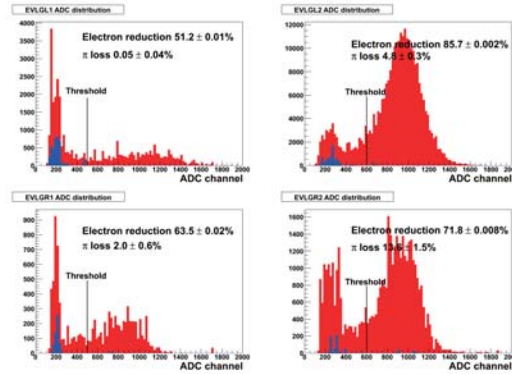


Fig.7 : EVLG2 ADC distribution. Blue fill : With π selection. Red fill : Without π selection. Threshold : setting value at experiment. N_{all} (with/without) : All number of EVLG ADC with / without π selection. $N_{\text{threshold}}$ (with/without) : Number above threshold. Statistics error Only.

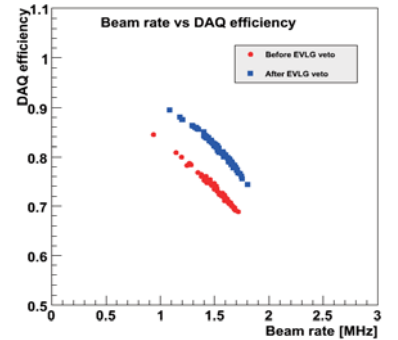


Fig.8 : Beam rate vs DAQ efficiency. Red circle : before EVLG veto. Blue cube : after EVLG veto.

✓ Summary

- Measurement of meson photoproduction with NKS2.
 - Elementary photoproduction on strangeness : $\gamma n \rightarrow K^0 \Lambda$.
 - A consistent cross section for $\gamma d \rightarrow \pi^+ \pi^- p n$ at wide energy range.
- We are improving NKS2.
 - Installation EVLG and VDC , and development tagger system.

- From EVLG ADC, number of e^+e^- background hitting this counter decreased to about 51% - 85%.
- Loss of π in one segment of EVLG is about 13% at the most.
- DAQ efficiency and beam intensity increased about 10% and 30% after incorporating EVLG as veto in the trigger, respectively.
- EVLG is effective in e^+e^- background reduction at trigger level.

Theory of Superconductivity in fullerides by the repulsive interaction model

Satoshi Yamazaki Yoshio Kuramoto
Department of Physics, Tohoku University

Introduction

A_3C_{60} $A=K, Rb, Cs$

- 3D crystal structure (bcc, fcc)
- Degenerate t_{1u} molecular orbitals (3 fold degeneracy)

Cs_3C_{60}

- Mott insulator at ambient pressure
- Superconducting (SC) state under applied pressure

The Coulomb repulsion plays a significant role.

Purpose

- Effect of the Coulomb repulsion for the SC state

→ Weak coupling approach

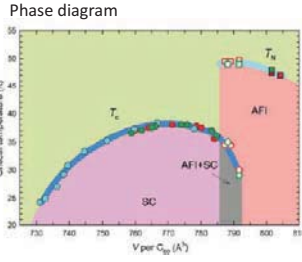
S. Rague, *et al*, Phys. Rev. B. **81** (2010) 224511.
S. Graser, *et al*, New J. Phys. **11**, 025016 2009.

- Effect of the crystal structure in A_3C_{60}

→ 3 dimensionality (bcc, fcc), orbital degeneracy

- Symmetry of superconductivity in A_3C_{60}

→ Point group T_h → A_g, E_g, T_g

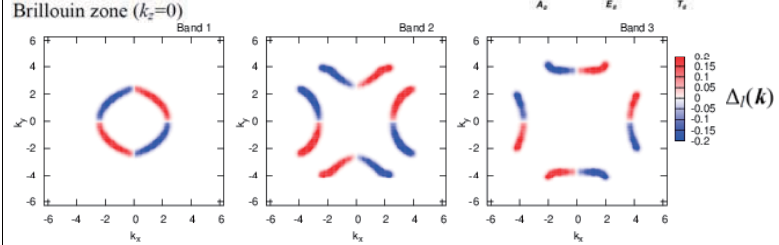


Y. Takabayashi, *et al*, Science **323**, 1585 (2009)

Result $U = U' = W = 0.5\text{eV}, J = K = 0$

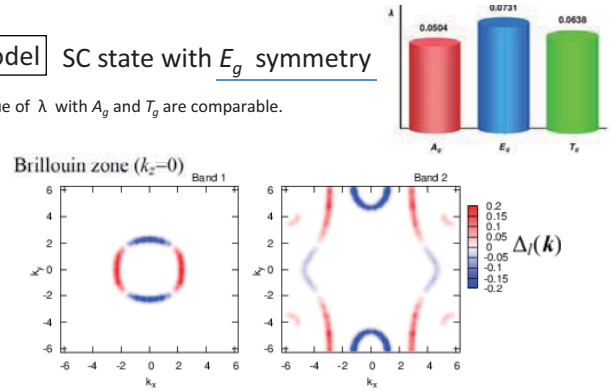
bcc model SC state with T_g symmetry

The SC state with A_g and E_g have the smaller value of λ .



fcc model SC state with E_g symmetry

The value of λ with A_g and T_g are comparable.



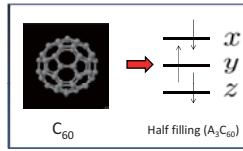
Model

Band width : $W=0.5\text{eV}$

Construct effective tight binding hamiltonian for t_{1u} orbitals

$$\mathcal{H}_0 = \sum_{\langle i,j \rangle} \sum_{\sigma, mn} t_{mn}(\alpha, \beta, \gamma) (c_{i\sigma mn}^\dagger c_{j\sigma mn} + h.c.)$$

$t_{mn}(\alpha, \beta, \gamma)$: hopping integral fitted by result of the first principle band calculation

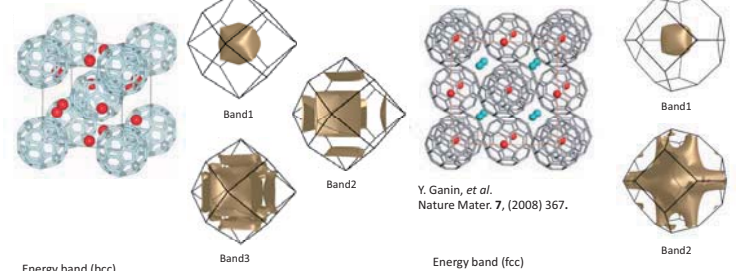


bcc structure

Fermi surface (bcc)

fcc structure

Fermi surface (fcc)



Y. Ganin, *et al*, Nature Mater. **7**, (2008) 367.

Kohn-Luttinger effective interaction (multi orbital)

W. Kohn and J.M. Luttinger, Phys. Rev. Lett. **15** (1965) 524.

$$V_{mn,m'n'}(k,k') = [\hat{U} + \hat{U}\hat{\chi}(k+k')\hat{U}]_{mn,m'n'}$$

$$U_{mm,mm} = U, U_{mn,mm} = U', U_{mm,mn} = J, U_{mn,mn} = K,$$

Intraorbital Interorbital Hund coupling Pair hopping

$$T_c \propto \exp(-1/\lambda), \lambda > 0$$

$$\lambda \Delta_l(k) = -\frac{1}{(2\pi)^3} \sum_{l'} \int_{\epsilon_{k'l'} = \epsilon_F} \frac{dS_{k'l'}}{v_{l'}(k')} V_{ll'}(k, k') \Delta_{l'}(k')$$

$$V_{ll'}(k, k') = \sum_{mn,m'n'} \bar{A}_l^m(k) \bar{A}_{l'}^{n'}(-k) V_{mn,m'n'}(k, k') A_{l'}^{m'}(k') A_l^m(-k').$$

Linear combination of the anisotropic SC state (O_h)

M. Sigrist and K. Ueda, Rev. Mod. Phys. **63**, 239 (1991).

T_{2g} symmetry

$$\omega = \exp(2\pi i/3)$$

	$\Delta_l(k)$	Degeneracy	Point group (irreducible representation)	
I	$k_y k_z + \omega k_x k_y + \omega^2 k_x k_z$	8	$D_{3d}(\Gamma_3^+)$	$d + id$ (chiral)
II	$k_y k_z + k_x k_y + k_x k_z$	4	$D_{3d}(\Gamma_1^+)$	
III	$k_y k_z$	3	$D_{4h}(\Gamma_4^+)$	d_{xy} (polar)
IV	$k_y k_z + i k_x k_z$	6	$D_{4h}(\Gamma_5^+)$	

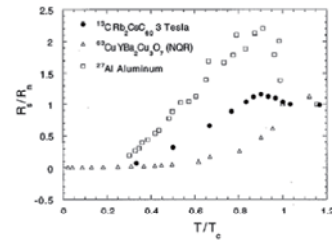
E_g symmetry

	$\Delta_l(k)$	Degeneracy	Point group (irreducible representation)
I	$k_x^2 - k_y^2$	3	$D_{4h}(\Gamma_3^+)$
II	$2k_z^2 - k_x^2 - k_y^2$	3	$D_{4h}(\Gamma_1^+)$
III	$k_x^2 + \omega k_y^2 + \omega^2 k_z^2$	2	$O_h(\Gamma_3^+)$

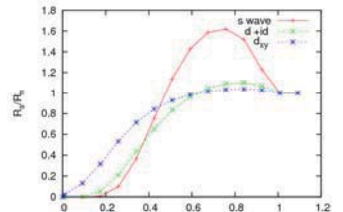
Temperature dependence of spin relaxation rate

$$R = \frac{1}{T_1 T} = \frac{1}{N} \sum_q |A_q|^2 \frac{\text{Im}[\chi(q, \omega_0, T)]}{\omega_0}$$

$$\Delta(T) = \Delta_0 \sqrt{1 - \frac{T}{T_c}} \quad T_c = \Delta_0 / 1.76$$



Experimental result



Calculation result (bcc model)

Conclusion & discussion

- Anisotropic SC state

bcc: T_g 3 fold degeneracy

fcc: E_g 2 fold degeneracy

- Degenerate SC state

→ This degeneracy is resolved below T_c .

- Small coherence peak caused by the $d + id$ (chiral) state (spin relaxation rate in nuclear magnetic resonance experimental)

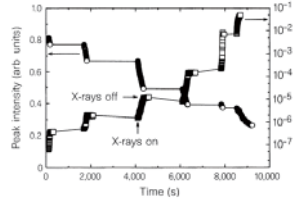


X-ray exposure effect on charge-orbital order in Fe-doped layered manganite $\text{La}_{0.5}\text{Sr}_{1.5}\text{Mn}_{1-x}\text{Fe}_x\text{O}_4$

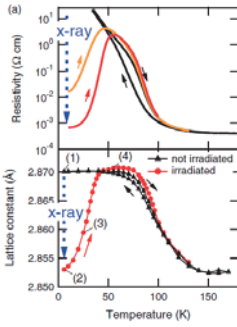
Y. Yamaki^{A, B}, Y. Yamasaki^B, H. Nakao^B, Y. Murakami^B, Y. Kaneko^{C, D}, Y. Tokura^{C, D, E, F}^A Department of science, Tohoku University, ^B Condensed Matter Research Center/Photon Factory, KEK, ^C Multiferroics Project, ERATO, ^D Cross-Correlated Materials Research Group, ^E Correlated Electron Research Group, RIKEN, ^F Department of Applied Physics, University of Tokyo

1. Introduction

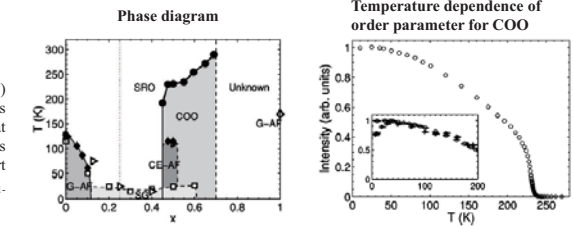
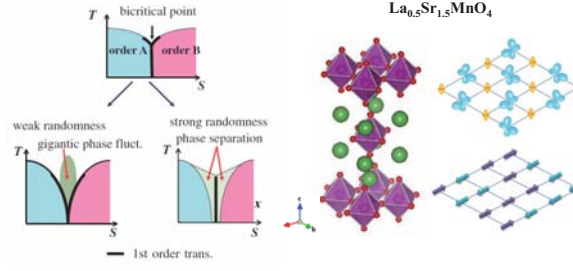
1-1. Photo induced phase transition

(1) $\text{Pr}_{0.7}\text{Ca}_{0.3}\text{MnO}_3$ 

- (1) Insulating antiferromagnetic state (charge ordering)
→ metallic ferromagnetic state
(2) Insulator → metal

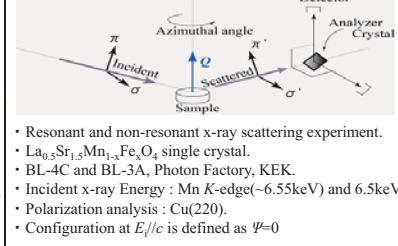
(2) W:VO_2 

1-2. Bicritical phase competition & object material

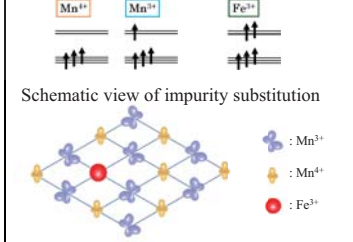


2. Experiments

Experimental configuration

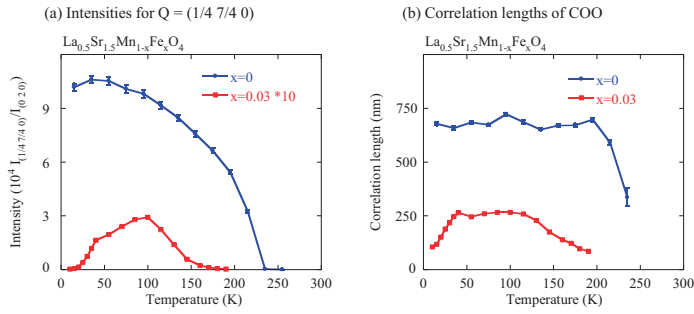


Spin configurations of Mn ions and Fe ion



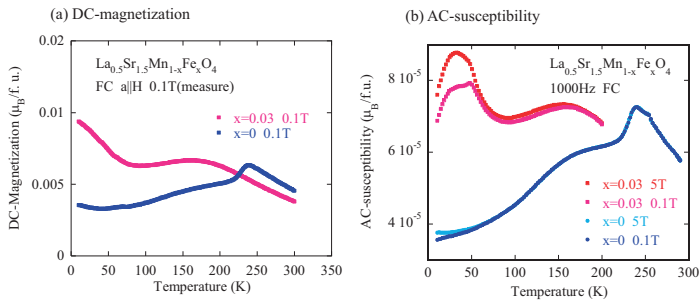
3. Results

3-1. Temperature dependence of the intensity for $Q = (1/4 \ 7/4 \ 0)$ and the correlation length of COO in $\text{La}_{0.5}\text{Sr}_{1.5}\text{Mn}_{1-x}\text{Fe}_x\text{O}_4$



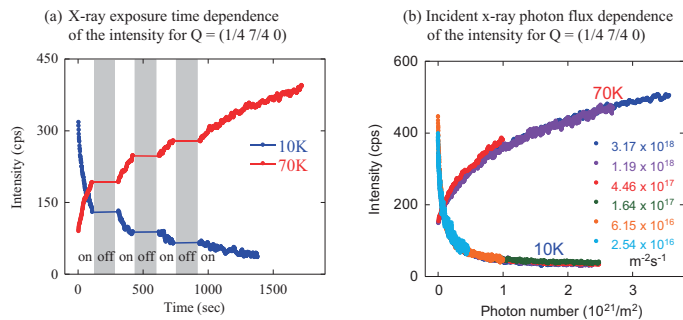
- (a) In Fe doped compound, the intensity of $Q = (1/4 \ 7/4 \ 0)$ which reflect lattice distortion accompanied with OO is about 30 times as small as that of pure compound.
(b) The correlation length of COO in Fe doped compound is about 3 times as short as that in pure compound.
→ COO state is strongly suppressed by substitution of Fe ions.

3-2. Temperature dependence of magnetic property in $\text{La}_{0.5}\text{Sr}_{1.5}\text{Mn}_{1-x}\text{Fe}_x\text{O}_4$



- (a) In Fe doped compound, DC-magnetization increases below 90 K.
(b) Temperature dependence of AC-susceptibility in Fe doped compound have peak structure around 40 K and this peak structure grows with increasing of magnetic field in field cooling.
→ Ferromagnetic component is emerged by substitution of Fe ions.

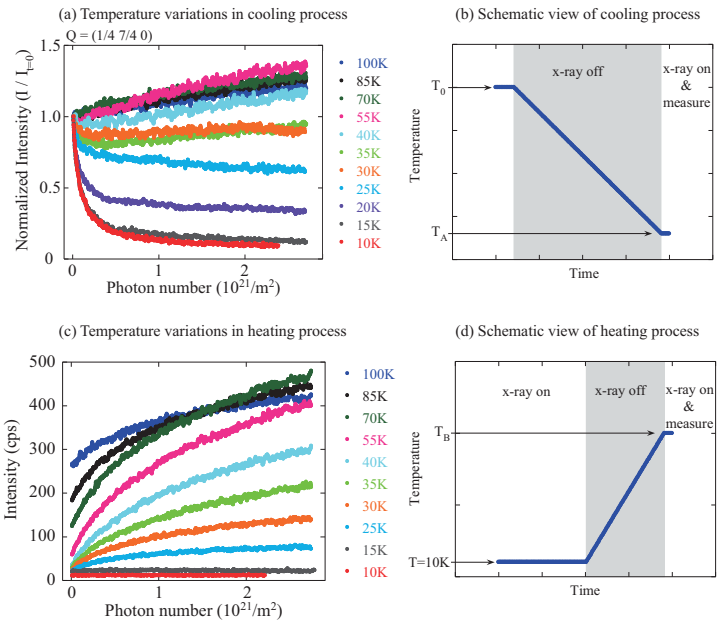
3-3. X-ray exposure effect on COO in $\text{La}_{0.5}\text{Sr}_{1.5}\text{Mn}_{0.97}\text{Fe}_{0.03}\text{O}_4$



Decreases (in 10 K) or increases (in 70 K) of intensities come about only when a sample is irradiated and changes of the intensities depend on only irradiated photon number.

→ This phenomenon is induced by x-ray exposure.

3-4. Photon number dependence of the intensity for $Q = (1/4 \ 7/4 \ 0)$ in $\text{La}_{0.5}\text{Sr}_{1.5}\text{Mn}_{0.97}\text{Fe}_{0.03}\text{O}_4$



- (a) The intensities are normalized by each intensity at $t = 0$. The intensities decrease below 25 K and increase above 55 K by x-ray exposure.
(c) The intensities all increase between 15 K and 100 K by x-ray exposure.
→ COO state is suppressed below 25 K and enhanced above 55 K by x-ray exposure.

4. Conclusion

In this study we have investigated photo induced phase transition in impurity doped manganite $\text{La}_{0.5}\text{Sr}_{1.5}\text{Mn}_{1-x}\text{Fe}_x\text{O}_4$. We assume that this photo induced phase transition would be caused by emerging of phase separation state induced by substitution of Fe ions.

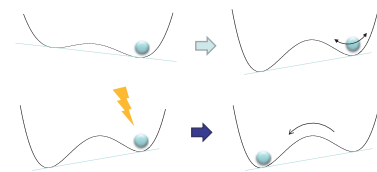
Cooling process

Intensity
10 K - 25 K
Monotone decreasing
30 K - 40 K
Decreasing, then increasing
55 K - 100 K
Monotone increasing

Heating process

Intensity
20 K - 100 K
Monotone increasing

Schematic view of photo induced phase transition in this system



Reference

- Saioa Cobo, Denis Ostrovskii, Sebastian Bonhommeau, Laure Vendier, Gabor Molnar, Lionel Salmon, Koichiro Tanaka, and Azzedine Bousseksou, J. AM. CHEM. SOC. **130**, 9019 (2008).
- V. Kiryukhin, D. Casa, J. P. Keimer, A. Vigliante, Y. Tomioka, and Y. Tokura, Nature **386**, 813 (1997).
- K. Shibuya, D. Okuyama, R. Kumai, Y. Yamasaki, H. Nakao, Y. Murakami, Y. Taguchi, T. Arima, M. Kawasaki, and Y. Tokura, Phys. Rev. B **84**, 165108 (2011).
- S. Larochelle, A. Mehta, L. Lu, P. K. Mang, O. P. Vajk, N. Kaneko, J. W. Lynn, L. Zhou, and M. Greven Phys. Rev. B **71**, 024435 (2005).
- Y. Tokura, JPSJ **75**, 011001 (2006)

Analysis of the hypernuclear γ -ray spectroscopy of $^{12}_{\Lambda}\text{C}$ and $^{11}_{\Lambda}\text{B}$ via the (π^+, K^+) reaction

K. Hosomi for KEK-E566 collaboration, Dept. of Physics, Tohoku Univ.

Introduction

ΛN interaction

Study of ΛN interaction is the first step to understand baryon-baryon interactions in a unified way beyond the well-known NN interaction. Instead of difficult ΛN scattering experiments due to a short lifetime of a Λ , it is a realistic and effective way to study ΛN interaction, in particular its spin-dependence, from precise structure of Λ hypernuclei.

Hypernuclear γ -ray spectroscopy

As shown in Fig. 1, the low-lying level structure of hypernuclei can be understood by combining a Λ in the 0s orbit and the other part called "core". Each state of the core nucleus (with non-zero spin) splits into two states in the hypernucleus. The spacing of a doublet splitting is determined by spin-dependent ΛN interactions and is in the order of a few 100 keV to a few 10 keV. Thus, high resolution γ -ray spectroscopy of hypernucleus with Ge detectors is a unique way to study spin-dependent ΛN interactions.

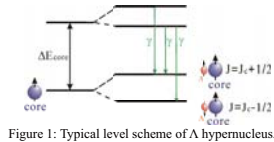


Figure 1: Typical level scheme of a hypernucleus.

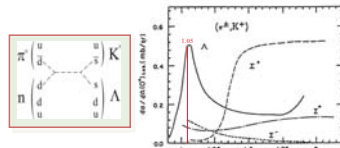
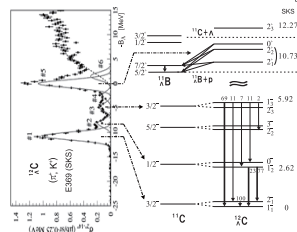


Figure 2: Diagram of $n(\pi^+, K^+)\Lambda$ and elementary cross section.

KEK-E566 experiment

The KEK-E566 experiment was performed at KEK-PS K6 beam line in 2005. In this experiment, (π^+, K^+) reaction was used on ^{12}C target to produce $^{12}_{\Lambda}\text{C}$ and $^{11}_{\Lambda}\text{B}$. The beam π^+ momentum was set to be 1.05 GeV/c because the elementary cross section has a sharp maximum at this energy as shown in Fig. 2. The outgoing K^+ had a momentum around 0.75 GeV/c. Fig.3 shows the expected level scheme of this experiment together with the spectrum taken in past experiment(E369).



Statistical analysis of human written language

Department of Physics , Sho Furuhashi

Introduction

Sentence length is a quantity showing the writing style. In Japanese, the distribution of sentence length (the number of characters) has some characteristic forms: **lognormal distribution**¹⁻³ and **Gamma distribution**².

There is no quantitative study about the features of the distribution of sentence length.

Purpose

We investigate the features of sentence length distribution by focusing on sentence structure.

- A) Lognormal distribution
- B) Gamma distribution

Conclusion

From the results,

- A) multiplicative process does not appear in dependency tree,
- B) and, sentence length is sum of exponential distributions.

Sentence length is similar to gamma distribution.

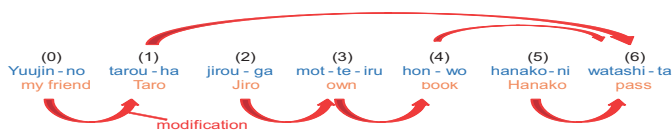
Reference

1. B. Yasumoto: Mathematical Linguistics No. 4 (1958) 20
2. K. Sasaki: Mathematical Linguistics No. 78 (1976) 13
3. H. Arai: Hitotsubashi University Journal 125 (2001) 205

Method and sample

Dependency tree

Sentence structure can be represented as tree based on dependency grammar between segments.



Sample

Aozora bunko: Online data base that contains public-domain texts.

(Total number of sentences: 116719)

Kyoto university text corpus: A text corpus compiled from the Mainichi Shinbun Newspaper.

(Total number of sentences: 38397)

Distribution of sentence length

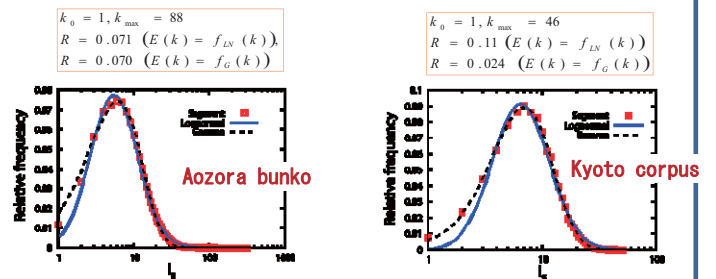
$$R = \sum_{k=k_0}^{k_{\max}} |O(k) - E(k)|$$

$O(k)$: Original data
 $E(k)$: Model distribution

$$f_{LN}(x) = \frac{1}{x\sqrt{2\pi\sigma^2}} \exp\left\{-\frac{(\ln x - \mu)^2}{2\sigma^2}\right\}$$

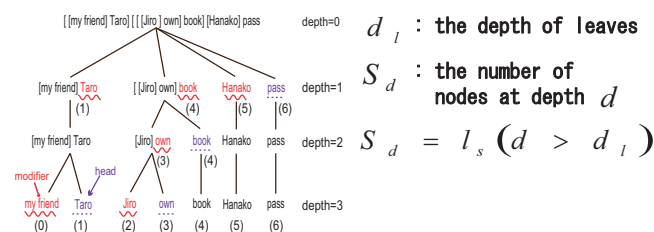
$$f_G(x) = x^{k-1} \frac{e^{-x/\theta}}{\Gamma(k)\theta^k}$$

Sentence length \rightarrow The number of segments l_s
 The distributions of $l_s \rightarrow$ **Gamma distribution**

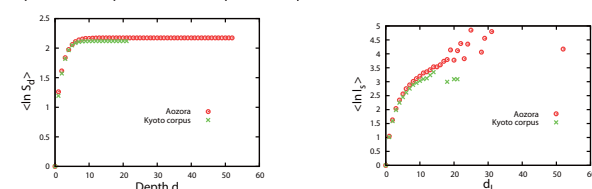


Result A

Multiplicative process $X_n = \prod_{i=0}^{n-1} \alpha_i X_0$ (α_i : positive random variable) is a model of **lognormal distribution**.

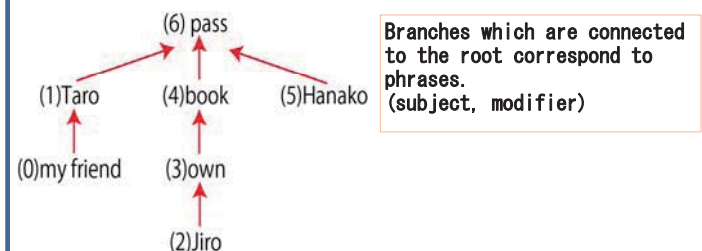


If the branching process is multiplicative process,
 $\langle \ln S_d \rangle \propto d$, $\langle \ln l_s \rangle \propto d_l$.

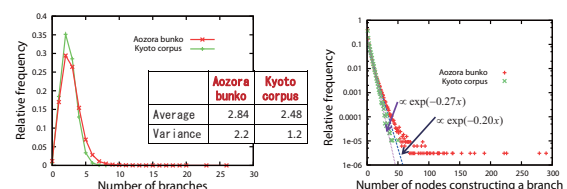


Result B

The sum of independent exponentially distributed random variables follows **gamma distribution**.



What forms are the distributions of the size and number of phrases?



Study of Λ photoproduction with Neutral Kaon Spectrometer 2

Takao Fujii for NKS2 Collaboration Department of Physics, Tohoku University

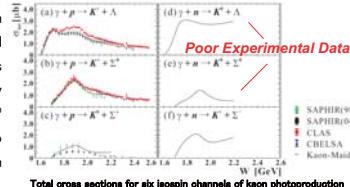
Introduction

Physics Motivation

Strangeness photoproduction in the threshold region plays the unique role for the investigation of hadron interactions and their structure. Among the six isospin channels of strangeness photoproduction, most of the experimental data have been obtained for the K^+ production channels. On the other hand, the experimental data for the reaction channels with K^0 was scarce because of their experimental difficulty (the right figure). So, the investigation of K^0 production channels is very important to understand the strangeness photoproduction mechanism.

Processes	Photon Energy Threshold
$\gamma + p \rightarrow K^+ + \Lambda$	911.1 MeV
$\gamma + p \rightarrow K^+ + \Sigma^0$	1046.2 MeV
$\gamma + p \rightarrow K^+ + \Sigma^+$	1047.5 MeV
$\gamma + n \rightarrow K^0 + \Lambda$	915.3 MeV
$\gamma + n \rightarrow K^0 + \Sigma^0$	1050.6 MeV
$\gamma + n \rightarrow K^0 + \Sigma^-$	1052.1 MeV

Photon energy threshold of six strangeness photoproduction channels



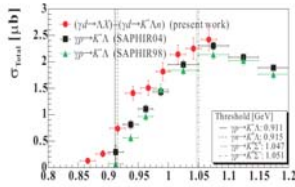
Total cross sections for six isospin channels of kaon photoproduction

Major decay mode of K^0 and Λ

$K^0 \rightarrow K_S^0$ (50%) / K_L^0 (50%)

K_S^0 ($c\tau = 2.68$ cm)	K_L^0 ($c\tau = 1534$ cm)
$K_S^0 \rightarrow \pi^+ \pi^-$ (69.2%)	$K_L^0 \rightarrow \pi^+ \pi^- \pi^0$ (40.6%)
$\rightarrow \pi^0 \pi^0$ (30.7%)	$\rightarrow \pi^+ \pi^- \pi^0$ (27.0%)
	$\rightarrow \pi^0 \pi^0 \pi^0$ (19.5%)
	$\rightarrow \pi^+ \pi^- \pi^0$ (12.5%)

$\Lambda \rightarrow p \pi^-$ (63.9%) / $n \pi^0$ (35.8%)



Photon Energy : E [GeV]
Total cross section of previous experiment (red point)

Previous experiment of $\gamma n \rightarrow K^0 \Lambda$ reaction

For the measurement of $\gamma n \rightarrow K^0 \Lambda$ reaction, we built the Neutral Kaon Spectrometer (NKS) and the NKS2 and have performed the photoproduction experiments with them. In the NKS and NKS2 experiment, the K^0 is detected as the K^0 's via its charged decay mode, and the Λ is also. Because no "pure" neutron target exists, the liquid Deuterium is used as the target.

The left figure is the total cross section of this reaction estimated from previous NKS2 experiment. It agrees with the total cross section of $\gamma p \rightarrow K^+ \Lambda$ reaction in this energy region ($0.8 < E_\gamma < 1.1$ GeV).

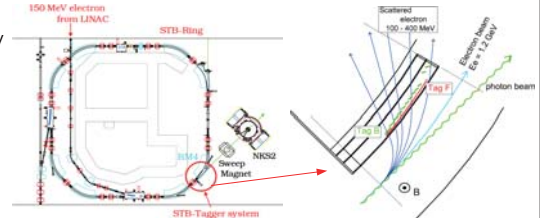
After this experiment, the central detectors of NKS2 has been upgraded to obtain the high statistical data.

Neutral Kaon Spectrometer 2 (NKS2)

Tagged photon beam

The experiment is performed using the real photon beam of the Research Center for Electron Photon Science (ELPH). Electron is accelerated to 1.2 GeV by the LINAC and the electron synchrotron called as STB-Ring. The photon is generated via bremsstrahlung, and scattered electron is tagged by the STB-Tagger system.

- Maximum energy of electron beam : 1.2 GeV
- Radiator : carbon fiber ($\phi 11 \mu\text{m}$)
- Tagger
 - TagF : 50 plastic scintillators
 - TagB : 12 plastic scintillators
- covered energy range : 6 MeV / 1 TagF counter
- Energy region : 0.8–1.1 GeV (Ee = 1.2 GeV)



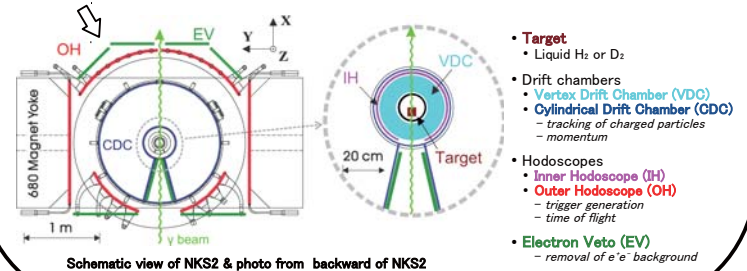
Schematic view of STB-Ring at ELPH and STB-Tagger system

NKS2

NKS2 is the one of the electromagnetic spectrometer designed for the measurement of $\gamma n \rightarrow K^0 \Lambda$ reaction. The Λ and the neutral kaon (K^0 's) are detected via their charged decay mode ($\Lambda \rightarrow p \pi^-$, $K_S^0 \rightarrow \pi^+ \pi^-$).

NKS2 consists of a large dipole magnet, two drift chambers, plastic scintillation hodoscopes and electron veto scintillation counters.

Liquid Hydrogen and Deuterium are used as the target.



Schematic view of NKS2 & photo from backward of NKS2

Experiment & Analysis

Experiment

The measurement of K^+ / K^0 and Λ photoproduction cross section with NKS2 has been performed using both of liquid Hydrogen and Deuterium target. The numbers of photons counted by Tagger were 7.7×10^{11} for the Hydrogen target period, and 8.8×10^{11} for the Deuterium target period respectively.

	Liq. H ₂	Liq. D ₂
Number of γ	7.7×10^{11}	8.8×10^{11}

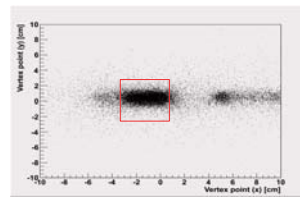
Calibration of Hodoscopes & Drift Chambers

Before the analysis of experimental data, the calibrations of hodoscopes and drift chambers were carried out. The time standard and the conversion factors (from TDC channels to time, from ADC channels to energy) were determined.

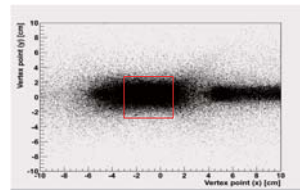
In the current analysis status, calibration of hodoscope is not finished completely. So, only forward region of NKS2, IH 1–5 ($\pm 54^\circ$ from beam direction) and OHV1–8 ($\pm 59^\circ$), are used in this analysis.

Vertex distribution

Here shows the vertex distribution of the horizontal plane for each target period. The target region is showed as the red square.



Vertex distribution of H2 target period



Vertex distribution of D2 target period

Particle identification (PID)

The particle identification is determined by the mass and the charge of the particle. The mass (m) is calculated by the correlation between the momentum (p) and velocity (β) as follows :

$$m^2 = p^2 \left(\frac{1}{\beta^2} - 1 \right)$$

The sign of charge is determined by the bending direction in the magnetic field.

- Definition of Pion

$$\beta^{-1} > 0.5$$

$$|p| > \frac{0.144}{\beta^{-1} - 0.2} - 0.08 \quad (0.5 < \beta^{-1} < 2.0)$$

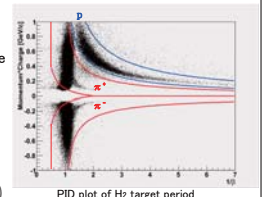
$$-0.5 < m^2 < 0.25$$

- Definition of Proton

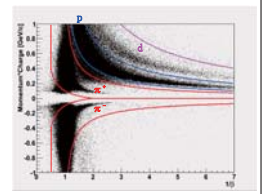
$$0.5 < m^2 < 1.8 \quad [(\text{GeV}/c^2)^2]$$

- Definition of Deuteron

$$1.8 < m^2 < 5.5 \quad [(\text{GeV}/c^2)^2]$$



PID plot of H2 target period



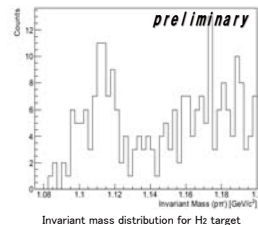
PID plot of D2 target period

Detection of Λ

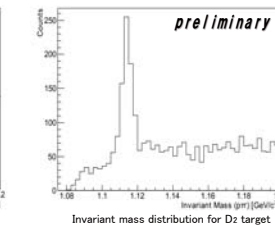
The Λ particle is detected as the peak in the $p \pi^-$ invariant mass spectrum. The invariant mass of $p \pi^-$ events is calculated as follows,

$$m_{p\pi^-}^2 = (m_p^2 + p_p^2 + m_{\pi^-}^2 + p_{\pi^-}^2) - |p_p + p_{\pi^-}|^2$$

Where m_p and m_{π^-} are the masses of the charged pion and proton from the PDG value, and p_p and p_{π^-} are the momenta of each particle. Here shows the $p \pi^-$ invariant mass distributions of the hydrogen and deuterium target period. In the H₂ data, the enhancement is appeared around 1.115 GeV/c².



Invariant mass distribution for H2 target



Invariant mass distribution for D2 target

In the D₂ data, a clear peak is obtained. The peak position is 1.114 GeV/c², and sigma of peak is 3.6 ± 0.2 MeV/c².

Summary

- Strangeness photoproduction, especially $\gamma n \rightarrow K^0 \Lambda$ reaction, plays the important role for the investigation of hadron interactions and their structure.
- NKS2 is the one of the electromagnetic spectrometer designed for the measurement of this reaction.
- Experiment has been performed using liquid H₂ and D₂ target.
- Calibration of data is performed, and the Λ particle is detected as the peak in the $p \pi^-$ invariant mass spectrum.
- Analysis for the calculation of the production cross section is being continued.

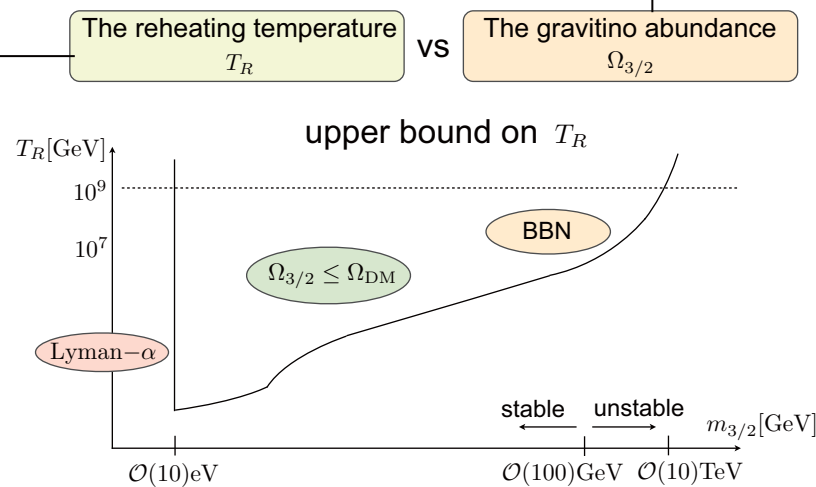
Cosmologically viable gauge mediation

Hiraku Fukushima (Physics, D1)

This work is based on the collaboration with Ryuichiro Kitano and Fuminobu Takahashi (Tohoku Univ.)

• Gravitino problem

: tension between



The gravitino abundance is constrained by the Big-Bang Nucleosynthesis (BBN) or observed Dark Matter density, which makes a severe upper bound on the reheating temperature. Gravitino with mass $O(10)\text{eV} \lesssim m_{3/2} \lesssim O(10)\text{TeV}$ is forbidden if we require $T_R \gtrsim 10^9\text{GeV}$.

• Constraint on $\Omega_{3/2}$

Gravitinos are produced in a thermal plasma after inflation, and their abundance is proportional to the reheating temperature.

$$\Omega_{3/2} \propto T_R.$$

Their abundance is constrained by

1. BBN (for unstable gravitinos);
Their decay should not destroy the success of the BBN!
Kawasaki, Kohri, Moroi(2005)
2. DM (for stable gravitinos);
Their abundance should not exceed the observed Dark Matter density today!
 $\Omega_{3/2} \lesssim \Omega_{DM} \simeq 0.2$.
Moroi, Murayama, Yamaguchi(1993)

• Why is a high reheating temperature needed?

In order to produce observed amount of baryon asymmetry by thermal leptogenesis, the mass of the right handed (s)neutrino must be larger than $10^9 - 10^{10}\text{GeV}$, which sets an upper bound on the reheating temperature, $T_R \gtrsim 10^9 - 10^{10}\text{GeV}$.

Davidson, Ibarra(2002), Giudice et al.(2004), Buchmuller, Bari, Plumacher(2004)

• Our Work

We found a novel scenario of gauge mediation:

- The reheating temperature can be raised as high as 10^{12}GeV with a gravitino mass $m_{3/2} \sim O(1)\text{GeV}$.
- Observed dark matter density $\Omega_{DM} \simeq 0.2$ is explained by thermally produced gravitino.
- Supersymmetry breaking vacuum is successfully selected as the temperature of the universe drops.

• Gravitino thermal production

The point:

Thermal production of the gravitino is highly suppressed if the messenger fields are thermalized.

K. Choi et al.(1999)

T_R can be raised as high as 10^{12}GeV without the over production of the gravitino!

	$T < M_{\text{mess}}$	$T > M_{\text{mess}}$
most effective process		
gravitino production rate	$\langle \sigma v_{\text{rel}} \rangle n_{\text{rad}} \sim \alpha_3 \frac{m_g^2}{m_{3/2}^2 M_{\text{pl}}^2} T^3$	$\langle \sigma v_{\text{rel}} \rangle n_{\text{rad}} \sim \alpha_3 \lambda^2 T$

• Cosmological evolution of pseudo-moduli

The model:

$$K = f^\dagger f + \bar{f}^\dagger \bar{f} + S^\dagger S - \frac{(S^\dagger S)^2}{\Lambda^2} + \dots$$

$$W = m^2 S - \lambda S f \bar{f} + c$$

S : SUSY breaking pseudo-moduli field

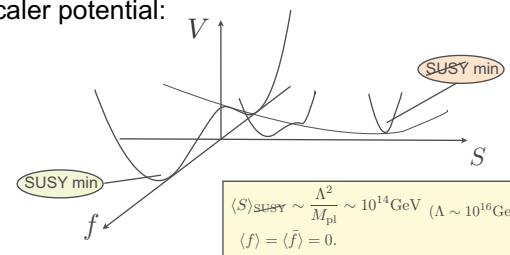
f, \bar{f} : messenger field

Λ : cut-off scale

m : SUSY breaking scale

c : supergravity correction

Scalar potential:



Thermal potential: $V_{\text{thermal}} \sim \lambda^2 T^2 |S|^2 + g^2 T^2 |f|^2$. g : SM gauge coupling

→ stabilize S, f at the origin.

For relatively small value of λ , S reaches SUSY minimum as the temperature of the universe drops.

Dalianis, Lalak(2011)

$$\Omega_{3/2} \propto \min[T_R, M_{\text{mess}}] \quad M_{\text{mess}}: \text{messenger mass scale}$$

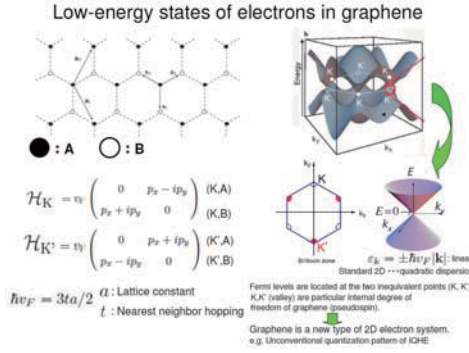
: independent of the reheating temperature for $T_R > M_{\text{mess}}$!

(different from the usual estimate $\Omega_{3/2} \propto T_R$.)

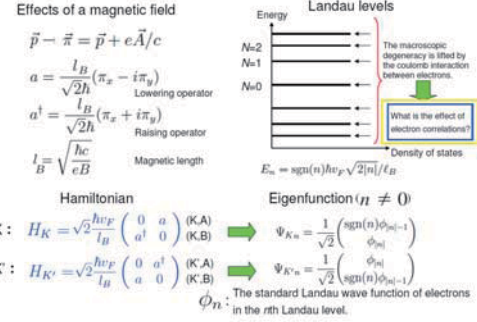
Using this formula, observed dark matter density is explained by thermally produced gravitino for $\lambda \sim 10^{-7}$ and $m_{3/2} \sim O(1)\text{GeV}$.

DMRG study of the ground state phase diagram of interacting massless Dirac fermions in graphene under magnetic field

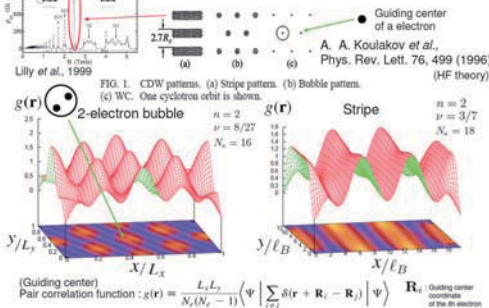
Tatsuya Higashi
Dept. of Physics, Tohoku U.



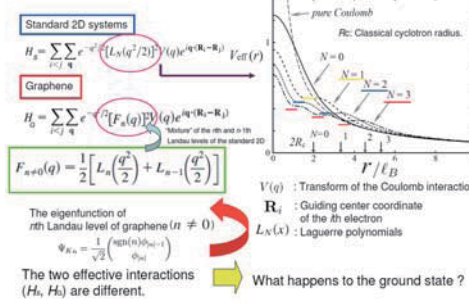
Electronic states of graphene in magnetic fields



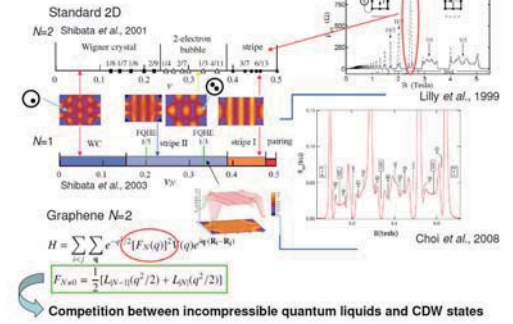
CDW ground states of standard 2D electron systems in high Landau levels



Effective inter-electron interaction of Mth Landau level



Motivation



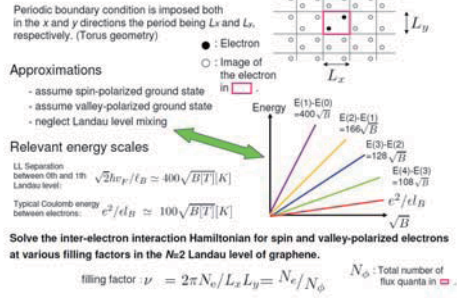
Purpose

Previous works

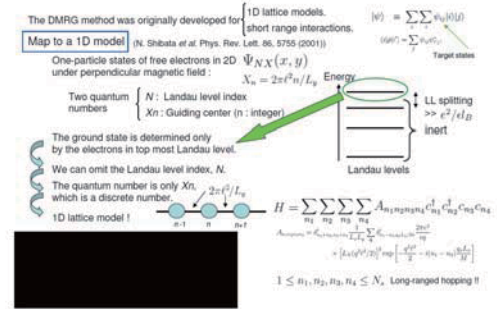
- Hartree-Fock theory (C.-H. Zhang et al., PRB 75, 245414 (2007))
 - Charge density waves (CDW) called stripes and bubbles realize in the ground state in high Landau levels of graphene.
 - But, quantum fluctuations are neglected in HF theory.
 - Can the CDW ground states survive under quantum fluctuations?
- Exact diagonalization (ED) (H. Wang et al., PRL 100, 116802 (2008))
 - The system size treated by ED is quite small.
 - The filling factor studied is limited.

In the present study, we investigate the ground state of graphene in the Landau levels of $N=2$ at various filling factors by the use of the DMRG method, and determine the ground state phase diagram.

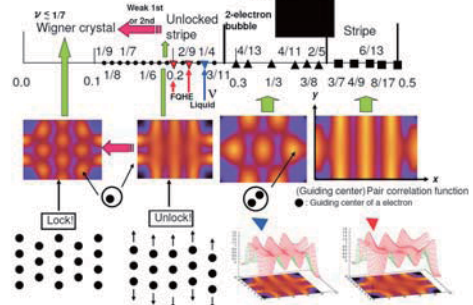
Model



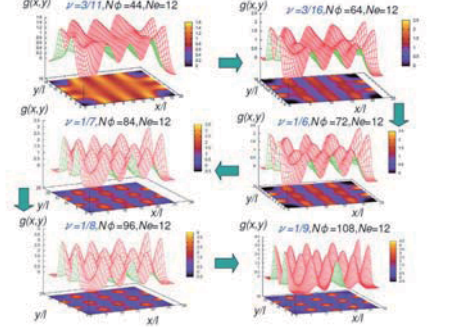
DMRG method for quantum 2D electron systems in a high magnetic field



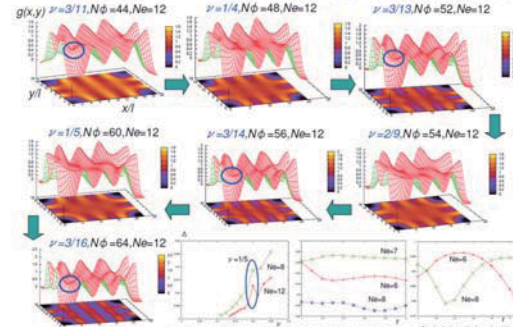
Ground state phase diagram of $N=2$ Landau level of graphene



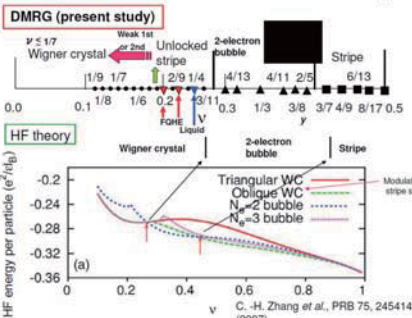
Ground state at low fillings



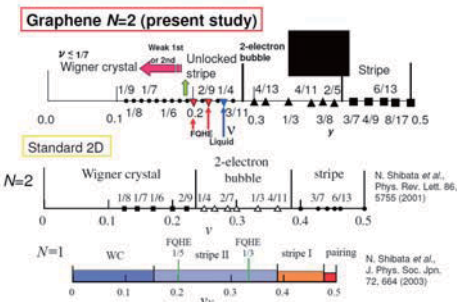
Ground state at low fillings



Comparison with the result of the HF theory



Comparison with standard 2D systems



Summary

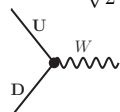
- By the use of DMRG method, we determined the ground state of electrons at various fillings of the $N=2$ Landau level of graphene.
- By analyzing the (guiding center) pair correlation function, we obtained the reliable ground state phase diagram.
- Possibility of realizing an unlocked stripe phase.
- Possibility of realizing a liquid ground state at $\nu=1/4$
- Possibility of realizing an incompressible quantum liquid state at $\nu=1/5$ and $2/9$

Study of $B^0 \rightarrow DK^{*0}(892)$ following by $D \rightarrow K^+ \pi^-$ at Belle

1. Motivation

Lagrangian of charged current weak interaction

$$\mathcal{L}_{int} = -\frac{g}{\sqrt{2}} (\bar{U}_L \gamma_\mu V_{CKM} D_L W_\mu^+) + h.c.$$



$U = (u, c, t)$

$D = (d, s, b)$

U_L, D_L : Left handed

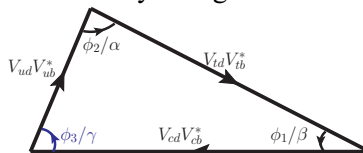
CKM (Cabbibo-Kobayashi-Masukawa) matrix

$$V_{CKM} = \begin{pmatrix} V_{ud} & V_{us} & V_{ub} \\ V_{cd} & V_{cs} & V_{cb} \\ V_{td} & V_{ts} & V_{tb} \end{pmatrix}$$

Unitarity condition

$$V_{CKM} V_{CKM}^\dagger = 1 \quad \text{row, 3 column elements} \quad V_{ud}V_{ub}^* + V_{cd}V_{cb}^* + V_{td}V_{tb}^* = 0$$

Unitarity triangle



$$\phi_1 = (21.15^{+0.90}_{-0.88})^\circ$$

$$\phi_2 = (89.0^{+4.4}_{-4.2})^\circ$$

$$\phi_3 = (68^{+13}_{-14})^\circ$$

Unitarity triangle is described on complex plane, and represents CP-violation.

To understand CP-violation, the angles of this triangle should be measured precisely.

In the present limits, measurement accuracy of ϕ_3 is not so good.

Need to study more for ϕ_3 .

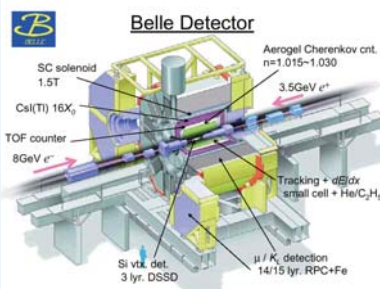


2. Belle



KEK@Tukuba

- KEKB-factory is a B facility.
- High energy electrons and positrons collide, and annihilate in pairs.
→ The pair annihilation produces a great deal of energy, and B particles are generated from the energy.
- There is the data of 1014 fb⁻¹ which is the largest in the world.
→ Good environment for B physics.



- Belle detector is to search the decays of B particles.
- Belle detector consists of many sub-detectors, and determines the particle type, momentum, charge, and so on.

3. Analysis

$$\phi_3 \equiv \arg \left(\frac{V_{ud}V_{ub}^*}{-V_{cd}V_{cb}^*} \right)$$

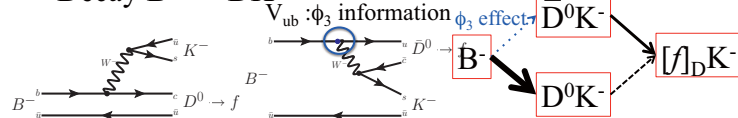
$$\sim -\arg(V_{ub})$$

ϕ_3 is measured with the decay include b→u transition.

e.c. $B^\pm \rightarrow D^{*0} K^\pm$

Influence of CP violation is expected to appear due to the interference between the two amplitudes of \bar{D}^0 and D^0 decays into a common final state.

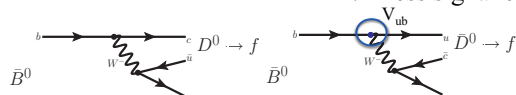
Decay $B^\pm \rightarrow DK^\pm$



Neutral B Decay $B^0 \rightarrow DK^{*0}$

1. More effect of ϕ_3

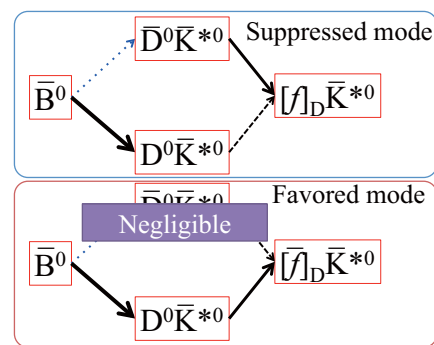
2. Less signal events → Large backgrounds



Where, I define $f = K^+ \pi^-$.

$$R_{DK^*} \approx \frac{\Gamma(B^0 \rightarrow [K^+ \pi^-]_D K^{*0}) + \Gamma(\bar{B}^0 \rightarrow [K^- \pi^+]_D \bar{K}^{*0})}{\Gamma(B^0 \rightarrow [K^- \pi^+]_D K^{*0}) + \Gamma(\bar{B}^0 \rightarrow [K^+ \pi^-]_D \bar{K}^{*0})}$$

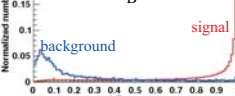
$$= r_S^2 + r_D^2 + 2kr_{SRD} \cos(\delta_S + \delta_D) \cos \phi_3$$



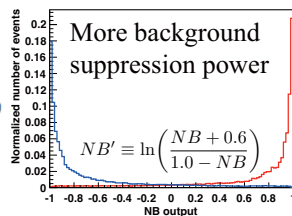
Background Subtraction

Common used parameter

1. Event topology
2. $\cos \theta_B$



New 7 parameters



Multivariate analysis is performed using nine variables. In order to optimize an acceptance and rejection, we use the NeuroBayes neural net package.

Since NB has a too peaky distribution to be fitted with simple analytic functions it is transformed to NB' to obtain a smoothed Gaussian like distribution.

4. Summary & Plan

I obtain R_{DK^*} upper limit.

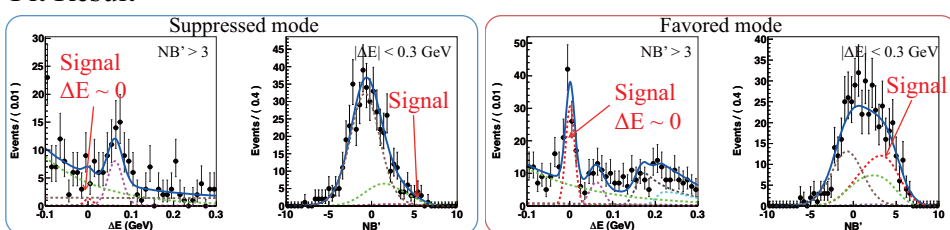
$$R_{DK^*} = (4.1^{+5.6+2.8}_{-5.0-1.8}) \times 10^{-2} < 0.16$$

I update R_{DK^*} upper limit world record. I'm writing the paper about this result.

Plan

- $B^0 \rightarrow [K_S \pi \pi]_D K^{*0}$ Dalitz plot analysis. ϕ_3 can be extracted with r_S, δ_S .
- I try to study about Belle vertex detector.

Fit Result



I perform 2D fit for ΔE and NB' .

In suppressed mode, there is no signal.

Study of light hypernuclei with the Stochastic Variational Method

Y. Nakagawa (Physics, Tohoku Univ.)

Light neutron-rich nuclei \leftrightarrow neutron halo structure of ${}^6\text{He}$, ${}^{11}\text{Li}$

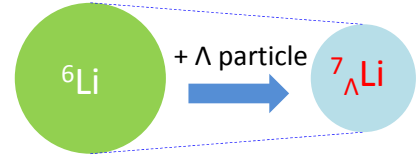
Λ hypernuclei \leftrightarrow Core shrinkage of ${}^7_\Lambda\text{Li}$

I. Tanihata,
J. Phys. G22, 157 (1996)
E. Hiyama et al., PRC59, 2351 (1999)
K. Tanida et al., PRL86, 1982 (2001)



Neutron-rich Λ hypernuclei

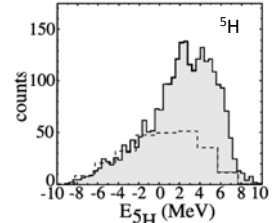
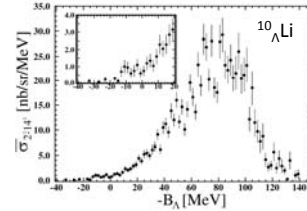
- ✓ Halo structure of nucleon + Λ particle ?
- ✓ Structure of core nucleus ?



✓ Production of neutron-rich hypernuclei

- ${}^{10}_\Lambda\text{Li}$ from KEK experiment
- ${}^6_\Lambda\text{H}$ from FINUDA experiment
- J-PARC experiment

P. K. Saha et al., PRL 94, 052502 (2005)
M. Agnello et al., PRL 108, 042501 (2012)
H. Tamura, JIMPA 24, 2101 (2009)



✓ Resonance state of super heavy hydrogen ${}^5\text{H}$ (N/Z = 4)

(Experiment) A. A. Korshennikov et al., PRL 87, 092501 (2001)

(Theory) t+n+n cluster model + Complex Scaling Method K. Arai, PRC 68, 034303 (2003)

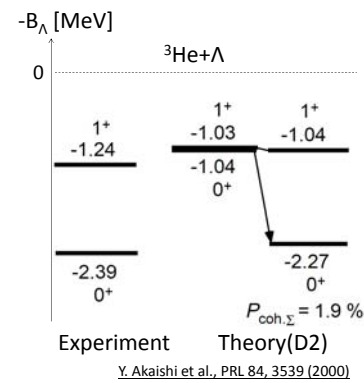
$\frac{1}{2}^+$ state \rightarrow (L,S) = (0,1/2), (1,1/2), (1,3/2), (2,3/2)

$E=1.7 \pm 0.3$ MeV, $\Gamma=1.9 \pm 0.4$ MeV

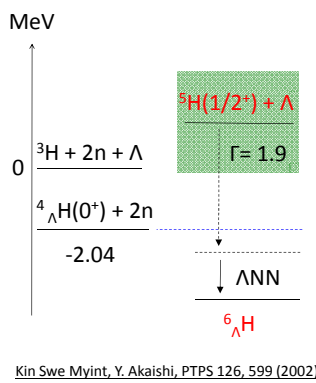
$E=1.59$ MeV, $\Gamma=2.48$ MeV

✓ $\Lambda\text{N}-\Sigma\text{N}$ coupling

- s-shell hypernuclei



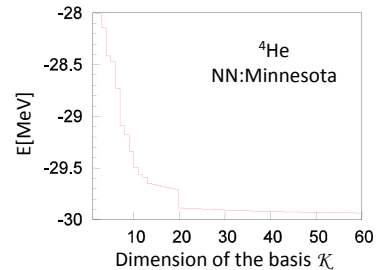
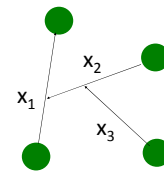
- n-rich hypernuclei



✓ Stochastic Variational Method

Y. Suzuki, K. Varga, LNP m54 (Springer, 1998)
Y. Suzuki et al., Few Body Syst. 42, 33 (2008)

1. Different \mathcal{N} sets of non-linear variational parameter are generated randomly.
2. By solving the \mathcal{N} eigenvalue problems, the corresponding energies are determined.
3. The parameter set producing the lowest energy is selected for \mathcal{K} th parameter set.
4. Increase \mathcal{K} to $\mathcal{K}+1$.
5. Return to 1.



✓ Variational trial function

$$\Psi_{JM_T T M_T} = \sum_{k=1}^{\mathcal{K}} c_k \psi_k \quad \psi_k = \mathcal{A} \left(\left[\phi_L^{space} \phi_S^{spin} \right]_{JM_T} \phi_{TM_T}^{isospin} \right)$$

$$\phi_{LM}^{space} = \exp \left(-\frac{1}{2} \sum_{i,j}^{A-1} A_{ij} \mathbf{x}_i \cdot \mathbf{x}_j \right) \theta_{LM}(\mathbf{x}) \quad \phi_{SM_S}^{spin} = \left| \left[\left[\frac{1}{2}, \frac{1}{2} \right]_{S_{12}}, \frac{1}{2} \right]_{S_{123}}, \dots \right]_{SM_S} \right\rangle$$

$$\theta_{LM}(\mathbf{x}) = [\mathcal{Y}_{L_1}(\mathbf{v}_1) \mathcal{Y}_{L_2}(\mathbf{v}_2)]_{LM} \quad \mathcal{Y}_{LM}(\mathbf{r}) \equiv r^L Y_{LM}(\hat{\mathbf{r}}) \quad \mathbf{v} = \sum_{i=1}^{A-1} u_i \mathbf{x}_i$$

\mathcal{A} : Antisymmetrizer

c_k : Linear variational parameters

A_{ij}, u_i : Non-linear parameters

✓ Application for s-shell hypernuclei

- Interaction

NN \rightarrow Minnesota, $\Lambda\text{N} \rightarrow$ Minnesota type
D. R. Thompson et al., NPA 286, 53 (1977)
H. Nemura et al., PTP 103, 929 (2000)

[MeV]	${}^3_\Lambda\text{H}$	${}^4_\Lambda\text{H}(0^+)$	${}^4_\Lambda\text{H}(1^+)$	${}^4_\Lambda\text{He}(0^+)$	${}^4_\Lambda\text{He}(1^+)$	${}^5_\Lambda\text{He}$
$B_\Lambda(\text{cal.})$	0.18	2.22	1.05	2.19	1.05	4.88
$B_\Lambda(\text{exp.})$	0.13	2.04	1.04	2.39	1.15	3.12

$u=1.5$ in $V_{\Lambda\text{N}}$

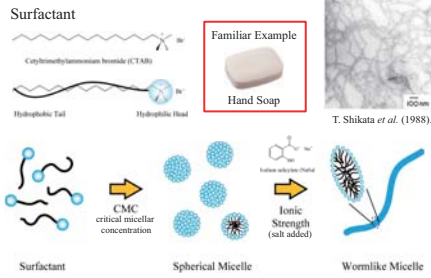
• Parameter "u" dependence in ΛN interaction

$$V = \left(V_R + \frac{1}{2}(1 + P_\sigma)V_t + \frac{1}{2}(1 - P_\sigma)V_s \right) \left(\frac{u}{2} + \frac{2-u}{2}P_r \right)$$

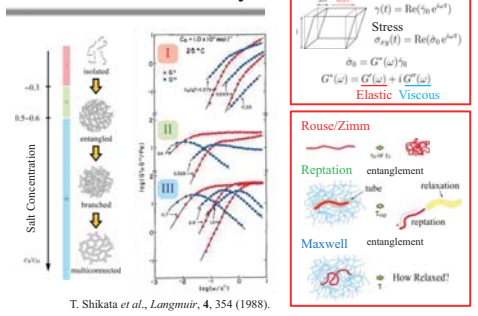
u	${}^3_\Lambda\text{H}$	${}^4_\Lambda\text{H}(0^+)$	${}^4_\Lambda\text{H}(1^+)$	${}^4_\Lambda\text{He}(0^+)$	${}^4_\Lambda\text{He}(1^+)$
1.5	0.18	2.22	1.05	2.19	1.05
2	0.20	2.28	1.14	2.24	1.12
2.5	0.22	2.37	1.22	2.35	1.19
(exp.)	0.13	2.04	1.04	2.39	1.15

$E_{\text{cal}}({}^2\text{H}) = -2.20$ MeV
 $E_{\text{cal}}({}^3\text{H}) = -8.38$ MeV
 $E_{\text{cal}}({}^3\text{He}) = -7.71$ MeV
 $E_{\text{cal}}({}^4\text{He}) = -29.94$ MeV

Wormlike Micelle ?



Linear Viscoelasticity

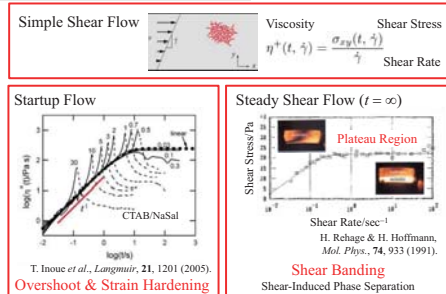


Statics & Dynamics of Wormlike Micellar Systems

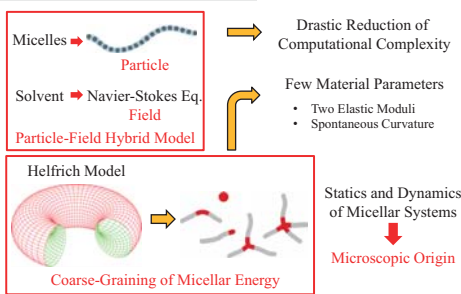
Masatoshi Toda
Physics, D3, Tohoku University

February 21, 2012

Nonlinear Viscoelasticity



Purpose of This Study



Particle-Field Hybrid Model

Micellar Component

$$\frac{d\mathbf{r}_i}{dt} = \mathbf{v}(\mathbf{r}_i) - \frac{\partial H(\{\mathbf{r}_j\})}{\partial \mathbf{r}_i} + \boldsymbol{\xi}_i$$

Solvent Component (Incompressible Linearized Navier-Stokes Eq.)

$$\frac{\partial \mathbf{v}}{\partial t} = -\nabla p + \eta \nabla^2 \mathbf{v} - \sum_{i=1}^N \frac{\partial H(\{\mathbf{r}_j\})}{\partial \mathbf{r}_i} \delta(\mathbf{r} - \mathbf{r}_i) + \boldsymbol{\theta}$$

Fluctuation-Dissipation Theorem

$$\langle \boldsymbol{\theta}(\mathbf{r}, t) \boldsymbol{\theta}(\mathbf{r}', t') \rangle = -2\eta T \nabla^2 \delta(\mathbf{r} - \mathbf{r}') \delta(t - t')$$

$$\langle \boldsymbol{\xi}_i(t) \boldsymbol{\xi}_j(t') \rangle = 2T \delta_{ij} \delta(t - t')$$

Y. Oono & K. F. Freed, *J. Chem. Phys.*, 75, 1009 (1981).

Coarse-Graining of Micellar Energy

Helfrich Elastic Energy

$$H_{\text{Bending}} = \int d\mathbf{a} \left[2\kappa(c_m(\mathbf{a}) - c_0)^2 + \bar{\kappa}c_G(\mathbf{a}) \right] = H_m + H_G$$

W. Helfrich, *Z. Naturforsch.*, 28c, 693 (1973).

Mean Curvature

$$c_m = \frac{1}{2} \left(\frac{1}{r_{\min}} + \frac{1}{r_{\max}} \right)$$

Gaussian Curvature

$$c_G = \frac{1}{r_{\min} r_{\max}}$$

Three Material Parameters

Cylindrical Bending Modulus

Saddle Splay Modulus

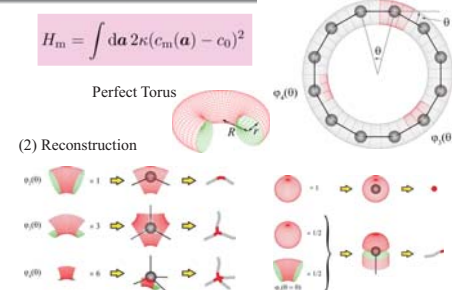
Spontaneous Curvature

$\mathbf{a} = (a_1, a_2)$

r_{\max}

r_{\min}

Mean Curvature Term



Gaussian Curvature Term

$H_G = \int d\mathbf{a} \bar{\kappa} c_G(\mathbf{a}) = 4\pi \bar{\kappa} (n_c - n_h) = 2\pi \bar{\kappa} \chi$

Elementary processes (Assumption)

(a) Scission $\Delta\chi = +2$

(b) Fusion $\Delta\chi = -2$

(c) End Interchange $\Delta\chi = 0$

(d) Bond Interchange $\Delta\chi = 0$

Topological Invariance

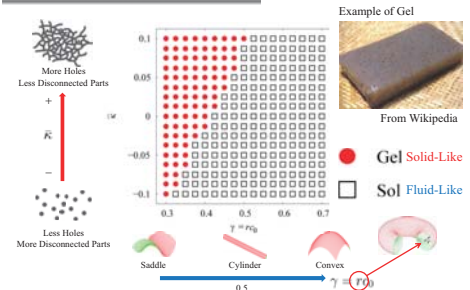
n_c : # of Disconnected Parts

n_h : # of Holes, i.e., Handles

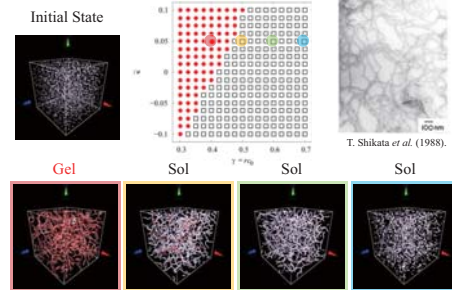
χ : Euler characteristics

M. Cates, *Macromolecules*, 20, 2289 (1987).

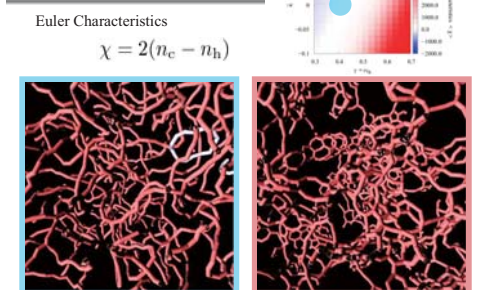
Phase Diagram



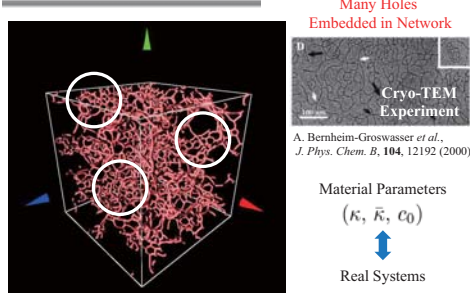
Internal Structures



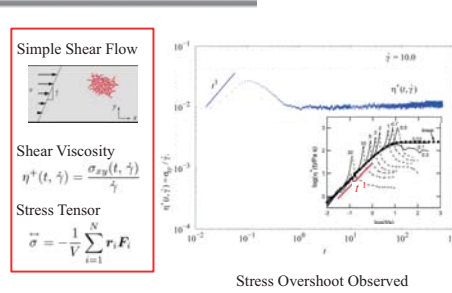
Topology of Networks



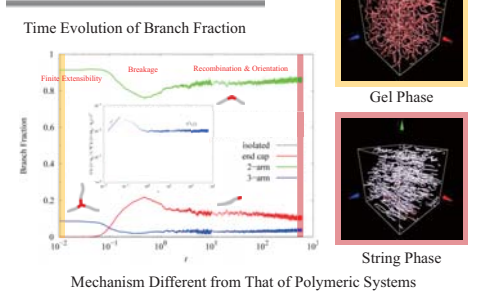
Comparison with Experiment



Startup Flow



Origin of Overshoot



Conclusions

Wormlike Micellar Systems

1. Statics
- Material Parameters $(\kappa, \bar{\kappa}, c_0)$ \Rightarrow Several Morphology (Sol/Gel)
2. Dynamics
- Overshoot \Rightarrow Mechanism Different from That of Polymer

Boundary state analysis on the equivalence of T-duality and Nahm transformation in superstring theory

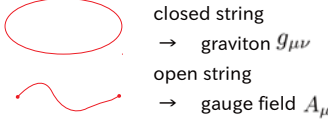
Y.Teshima

In collaboration with T. Asakawa, S. Wataura and U. Carow-Watamura
arXiv:1201.0125v2 [hep-th]

Introduction

String theory

A theory which regard elementary objects as a string



String theory is a candidate of quantum gravity
but defined only perturbatively

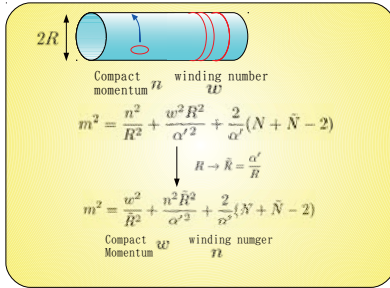
→ we want to find the stringy geometry

(cf. general relativity ↔ Riemannian geometry)

→ This must encode the property of T-duality

T-duality

Spectrum of the closed string



The theory with radius R
= The theory with radius \tilde{R}

→ T-duality

D-brane

the hypersurface where open strings end.

The direction along the D-brane

→ Neumann condition

The direction perpendicular to the D-brane

→ Dirichlet condition

p dimensional D-brane = Dp-brane

The effective theory on N D-branes

= $U(N)$ gauge theory



T-duality and D-brane

Dirichlet and Neumann condition are

exchanged under the T-duality

→ $\begin{cases} \text{T-dual normal to the D-brane} \\ \text{Dp} \rightarrow \text{D}(p+1) \\ \text{T-dual along the D-brane} \\ \text{Dp} \rightarrow \text{D}(p-1) \end{cases}$

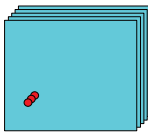
Bound state of D-branes

Bound state of the D-branes

Correspond to a soliton solution
in the effective theory

※ Topological charge

= the number of D-branes



Various approaches to discuss T-duality

Buscher rule

Nahm transformation

Boundary state analysis

Hori formula

→ We analyzed the compatibility among them.

Nahm transformation

Original Nahm transformation

$U(N)$ Gauge theory on T^4 → $U(k)$ Gauge theory on \tilde{T}^4
instanton instanton

D-brane interpretation

N D4 / k D0

k D4 / N D0

Consistent with T-duality

Generalization of Nahm transformation to 2d

We consider the classical solution on T^2

= constant magnetic flux

$$U(N), C_1 = \int \text{tr} \frac{F}{2\pi} = k$$

$$\Rightarrow A_1 = 0, A_2 = \frac{k}{2\pi N R_1 R_2} x_1, F_{12} = \frac{k}{2\pi N R_1 R_2}$$

There is a freedom to add constants \tilde{x}_1, \tilde{x}_2

$$\rightarrow A_1 = \tilde{x}_1, A_2 = \frac{k}{2\pi N R_1 R_2} x_1 + \tilde{x}_2$$

$$\tilde{x}_\mu \sim \tilde{x}_\mu + 1/R_\mu$$

$$\Rightarrow \tilde{x}_\mu \text{ is the coordinate on } \tilde{T}^2$$

We will construct the gauge field on \tilde{T}^2 by

Dirac zero mode ($\mathcal{D}_x \psi = 0$).

\mathcal{D}_x : The Dirac operator on $T^2 \times \tilde{T}^2$

$$\mathcal{D}_x = \begin{pmatrix} \partial_x & 0 \\ 0 & \partial_x \end{pmatrix}, \mathcal{D}_x = -i(\partial_1 - iA_1 - i\tilde{x}_1) - (\partial_2 - iA_2 - i\tilde{x}_2),$$

$$= \begin{pmatrix} \partial_x & 0 \\ 0 & \partial_x \end{pmatrix}, \mathcal{D}_x = -i(\partial_1 - iA_1 - i\tilde{x}_1) + (\partial_2 - iA_2 - i\tilde{x}_2)$$

$\mathcal{D}_x \psi = 0$ has k solutions

→ construct $N \times k$ matrix

$$\rightarrow \psi = (\psi_1 \cdots \psi_k) \cdot N$$

k

The gauge field on \tilde{T}^2 is

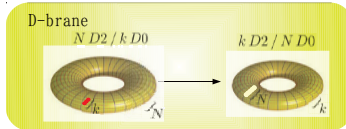
$$\tilde{A}_\mu(\tilde{x}) = i \int_{T^2} d^2 x \psi^\dagger \tilde{\partial}_\mu \psi$$

$$\rightarrow \tilde{A}_1 = 0, \tilde{A}_2 = -\frac{N}{2\pi k R_1 R_2} \tilde{x}_2, \tilde{F}_{12} = -\frac{N}{2\pi k R_1 R_2}$$

$$\rightarrow U(k), C_1 = -N$$

Atiyah-Singer family index theorem

$$\int_{T^2} \text{ch}(\mathcal{P}) \text{ch}(E) = \text{ch}(\tilde{E})$$

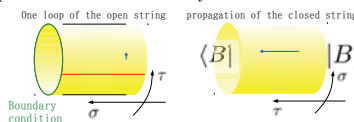


Nahm transformation:
 N D2 / k D0 → k D2 / N D0

Boundary state on the torus

Boundary state

partition function of cylinder



The initial, final state of the closed string
associated with

the boundary condition of the open string

D-brane

ex) In case of Dirichlet condition

$$\frac{\partial X^1}{\partial \tau} \Big|_{\sigma=0} = 0 \iff \frac{\partial X^1}{\partial \sigma} \Big|_{\tau=0} |B\rangle = 0$$

Calculate the boundary state with the flux associated
with $(U(N), C_1 = k)$ on T^2

$$|B\rangle = \sqrt{\frac{1}{2\pi\alpha'}} \exp \left\{ -\frac{1}{2\pi\alpha'} \sum_{n \neq 0} \frac{1}{n} \exp(-in\tau) \exp(in\sigma) \right\} \times \left[\begin{pmatrix} \alpha_0^1 & \alpha_0^2 \\ \alpha_1^1 & \alpha_1^2 \\ \alpha_2^1 & \alpha_2^2 \end{pmatrix} \begin{pmatrix} \alpha_0^1 & \alpha_0^2 \\ \alpha_1^1 & \alpha_1^2 \\ \alpha_2^1 & \alpha_2^2 \end{pmatrix} \right]$$

T-dual transformation $\begin{cases} m^M \rightarrow n^M \\ \alpha_0^M \rightarrow \alpha_0^M \\ \alpha_1^M \rightarrow \alpha_1^M \\ \alpha_2^M \rightarrow \alpha_2^M \end{cases}$

$$|B\rangle_{12} = \sqrt{N^2 + (\tilde{\alpha}_0^1 k)^2} \sum_{n \in \mathbb{Z}} \exp(-in\tau) \exp(in\sigma) \times \left[\begin{pmatrix} \alpha_0^1 & \alpha_0^2 \\ \alpha_1^1 & \alpha_1^2 \\ \alpha_2^1 & \alpha_2^2 \end{pmatrix} \begin{pmatrix} \alpha_0^1 & \alpha_0^2 \\ \alpha_1^1 & \alpha_1^2 \\ \alpha_2^1 & \alpha_2^2 \end{pmatrix} \right]$$

$$\rightarrow U(k), C_1 = -N$$

Consistent with the 2d Nahm transformation

Fermionic Part

NS-NS sector $(\psi_r, \tilde{\psi}_r)$ ($r \in \mathbb{Z} + \frac{1}{2}$)

R-R sector $(\psi_n, \tilde{\psi}_n)$ ($n \in \mathbb{Z}$)

We consider the R-R zero modes.
(Other modes contribute to the boundary state
in the same way to the bosonic sector.)

The state

$$\theta^\mu = (\psi_0^\mu + i\tilde{\psi}_0^\mu) / \sqrt{2}$$

$$|Dp\rangle = \theta_1^\dagger \theta_{p-1}^\dagger \cdots \theta_0^\dagger |C\rangle \quad \{\theta^\mu, \theta_\nu^\dagger\} = \delta_\nu^\mu$$

$|C\rangle$: the vacuum state

satisfies the boundary condition for a Dp brane.

$$\begin{aligned} \theta_\alpha^\dagger |Dp\rangle &= 0, \quad (\alpha = 0, \dots, p), \\ \theta_i |Dp\rangle &= 0, \quad (i = p+1, \dots, 9) \end{aligned}$$

Then the D2/D0 bound state is

$$\begin{aligned} |D2D0\rangle &= N e^{2\pi\alpha' F_{12} \theta^1 \theta^2} |D2\rangle \\ &= N(1 + F_{12} \theta^1 \theta^2) |D2\rangle \\ &= N |D2\rangle + k |D0\rangle \end{aligned}$$

T-duality transformation is represented by an operator:

$$\mathcal{T}_\alpha = \theta^\alpha - \theta^{\alpha\dagger}$$

Then the T-duality of $|D2D0\rangle$ is

$$|D2D0\rangle' = \mathcal{T}_2 \mathcal{T}_1 |D2D0\rangle = k |D2\rangle - N |D0\rangle$$

$$(N, k) \rightarrow (k, -N)$$

RR-charges of D-branes are measured by the coupling
to closed-string states of RR-potentials.

RR-state can be defined as

$$\langle A | = \langle [C] | A = \langle [C] | \sum A_{\mu_1 \cdots \mu_p} \theta^{\mu_1} \cdots \theta^{\mu_p}$$

And the T-duality transformation is

$$\begin{aligned} \langle A | \mathcal{T}_1^\dagger \mathcal{T}_2^\dagger &= \langle [C] | \left(A^{(0)} + A_1^{(1)} \theta^1 + A_2^{(1)} \theta^2 + A_{12}^{(2)} \theta^1 \theta^2 \right) \mathcal{T}_1^\dagger \mathcal{T}_2^\dagger \\ &= \langle [C] | \left(A^{(0)} \theta^1 \theta^2 - A_1^{(1)} \theta^2 + A_2^{(1)} \theta^1 + A_{12}^{(2)} \right) \end{aligned}$$

Then the T-duality rule for the RR-potentials is

$$A^{(0)} = -A_{12}^{(2)}, A_1^{(1)} = A_2^{(1)}, A_2^{(1)} = -A_1^{(1)}, A_{12}^{(2)} = A^{(0)}$$

The coupling of RR-potentials to D-brane is

$$I_{CS} = \langle A | [GSO] c_0 \tilde{c}_0 | B \rangle, \quad \begin{aligned} [GSO]: & \text{GSO projection} \\ |B\rangle: & \text{full boundary state} \\ c_0, \tilde{c}_0: & \text{ghost} \end{aligned}$$

We can prove that I_{CS} is invariant under T-duality:

$$\begin{aligned} I_{CS} &= \langle A | [GSO] \mathcal{T}_1^\dagger \mathcal{T}_2^\dagger c_0 \tilde{c}_0 | B \rangle \quad (1 = \mathcal{T}^\dagger \mathcal{T}, \mathcal{T} = \mathcal{T}_2 \mathcal{T}_1) \\ &= \langle A | \mathcal{T}^\dagger [GSO] c_0 \tilde{c}_0 \mathcal{T} | B \rangle \quad (\text{Since } \mathcal{T} \text{ commutes with } [GSO].) \\ &= \langle A' | [GSO] c_0 \tilde{c}_0 | B' \rangle \end{aligned}$$

I_{CS} is invariant under T-duality transformation.

Hori formula

I_{CS} can be written in terms of differential forms as

$$\begin{aligned} I_{CS} &= \mu_2 \int_M A \wedge \text{Tr}_N (e^{2\pi\alpha' F}) \\ &= \mu_2 \int_M (A^{(0)} + \cdots + A_{12}^{(2)} dx^1 \wedge dx^2) (N + k dx^1 \wedge dx^2) \\ &= N \mu_2 \int_M A_{12}^{(2)} dx^1 \wedge dx^2 + k \mu_0 \int_R A^{(0)} \end{aligned}$$

In terms of differential forms,

T-duality is represented by Hori formula:

$$A' = - \int_{T^2} A e^{dx^1 \wedge dx^2}$$

By using this, the Chern-Simons term is transformed as

$$\begin{aligned} \mu_2 \int_M A' \wedge \text{Tr}_k (e^{2\pi\alpha' \tilde{F}}) & \quad (\tilde{M} = \mathbb{R} \times \tilde{T}^2) \\ &= \mu_2 \int_{\tilde{M}} (A^{(0)} + \cdots + A_{12}^{(2)} dy^1 \wedge dy^2) (k - N dy^1 \wedge dy^2) \\ &= N \mu_2 \int_{\tilde{M}} A_{12}^{(2)} dy^1 \wedge dy^2 + k \mu_0 \int_R A^{(0)} \end{aligned}$$

→ invariant under T-duality

Summary and discussion

- Nahm transformation was extended naively to 2d.
- It consistent with T-duality in string theory.
(N, k) → ($k, -N$)
- It indicates \mathbb{Z}_4 -duality nature of T-duality
cf. Fourier transformation $(x, p) \rightarrow (p, -x) \rightarrow (-x, -p)$
- We introduced the T-duality operator which act both
on the boundary state and the RR q-form state.
- We clarified the relationship between T-duality rule at the
superstring level and that at the low energy effective theory

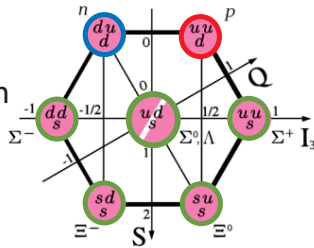
Description of single- Λ hypernuclei with relativistic point coupling model

Department of Physics, Nuclear Theory group
Yusuke Tanimura

Hypernuclear physics

Normal nuclei consist of **protons** and **neutrons**.

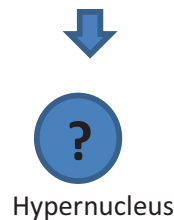
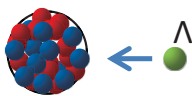
Nuclei containing baryons with strangeness (**hyperons**) are called hypernuclei



Λ **particle** may change properties of normal nuclei (*impurity effect*) such as

- Shape and radius
- Cluster structure
- Collective excitations
- ...

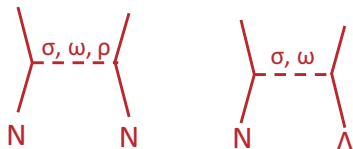
Normal nucleus



Response of nuclei to an addition of a Λ
→ **Structure information which cannot be seen with normal probes**

Relativistic mean field (RMF) calculations for hypernuclei

NN and $N\Lambda$ interactions are described as meson exchanges



No **relativistic 3-D mesh calculation** has been realized because of “variational collapse”,
Y. Zhang, et al, Int. J. Mod. Phys. E **19** 55 (2010)
K. Hagino and Y. T., PRC **82**, 057301 (2010)

→RMF is not developed as much as non-relativistic mean field calc., in which 3-D calculations are performed using **zero range interactions**.

→ **Relativistic zero range $N\Lambda$ interaction** may be useful in future relativistic calculation for hypernuclei

Model Lagrangian

We have proposed a new nucleon- Λ interaction:

Relativistic point coupling (RPC) model

$$\begin{aligned}\mathcal{L}_{\text{int}}^{N\Lambda} &= \mathcal{L}_{4f}^{N\Lambda} + \mathcal{L}_{\text{der}}^{N\Lambda} + \mathcal{L}_{\text{ten}}^{N\Lambda} \\ \mathcal{L}_{4f}^{N\Lambda} &= -\alpha_S^{(N\Lambda)} (\bar{\psi}_N \psi_N) (\bar{\psi}_\Lambda \psi_\Lambda) \\ &\quad -\alpha_V^{(N\Lambda)} (\bar{\psi}_N \gamma_\mu \psi_N) (\bar{\psi}_\Lambda \gamma^\mu \psi_\Lambda) \\ \mathcal{L}_{\text{der}}^{N\Lambda} &= -\delta_S^{(N\Lambda)} (\partial_\mu \bar{\psi}_N \psi_N) (\partial^\mu \bar{\psi}_\Lambda \psi_\Lambda) \\ &\quad -\delta_V^{(N\Lambda)} (\partial_\mu \bar{\psi}_N \gamma_\nu \psi_N) (\partial^\mu \bar{\psi}_\Lambda \gamma^\nu \psi_\Lambda) \\ \mathcal{L}_{\text{ten}}^{N\Lambda} &= -\alpha_T^{(N\Lambda)} (\bar{\psi}_N \sigma^{\mu\nu} \psi_N) (\partial_\mu \bar{\psi}_\Lambda \gamma_\nu \psi_\Lambda)\end{aligned}$$



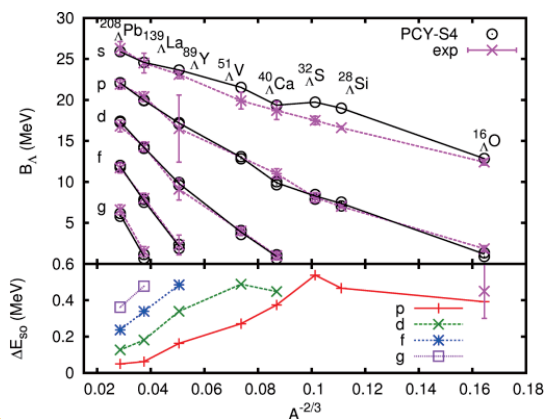
- An extension of RPC for NN interaction (Buervenich et al., PRC65, 044308('02))
- Finite range meson exchange is simulated by contact couplings between the fermions
- Zero range is numerically easy to handle**

Results

5 parameters are fitted to the experimental data of Λ binding energies of $^{16}_\Lambda\text{O}$, $^{40}_\Lambda\text{Ca}$, $^{51}_\Lambda\text{V}$, $^{89}_\Lambda\text{Y}$, $^{138}_\Lambda\text{La}$, $^{208}_\Lambda\text{Pb}$.

Comparison between theoretical and experimental Λ binding energies

(Y. T. and K. Hagino, Phys. Rev. C **85**, 014306 (2012))



- Well reproduces the experimental data ($\chi^2_{\text{dof}} = 0.92$)
- Suitable for 3D calculations or beyond-mean-field methods because of its numerical simplicity

“Study of $B^\pm \rightarrow DK^\pm$, $D \rightarrow K_S K^\pm \pi^\mp$ for the measurement of CP-violating angle ϕ_3 , and $D^{*\pm} \rightarrow D\pi^\pm$, $D \rightarrow K_S K^\pm \pi^\mp$ for the modeling of $D \rightarrow K_S K^\pm \pi^\mp$ Dalitz plane



1. Motivation & Theory

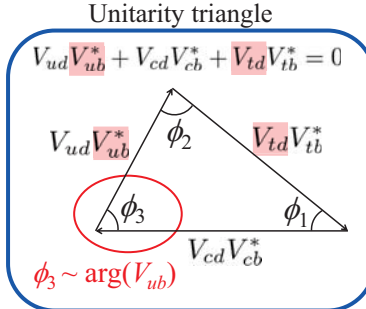
CKM (Cabibbo-Kobayashi-Maskawa) Matrix

$$V = \begin{pmatrix} V_{ud} & V_{us} & V_{ub} \\ V_{cd} & V_{cs} & V_{cb} \\ V_{td} & V_{ts} & V_{tb} \end{pmatrix}$$

$$= \begin{pmatrix} 1 - \frac{\lambda^2}{2} & \lambda & A\lambda^3(\rho - i\eta) \\ -\lambda & 1 - \frac{\lambda^2}{2} & A\lambda^2 \\ A\lambda^3(1 - \rho - i\eta) & -A\lambda^2 & 1 \end{pmatrix} + \mathcal{O}(\lambda^4)$$

$\lambda = \sin \theta_c \sim 0.22$

Unitarity
 $VV^\dagger = 1$



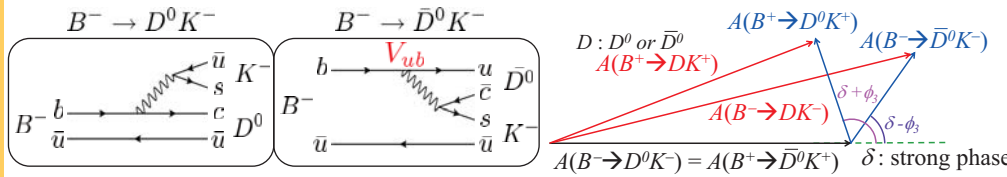
Unitarity triangle is described on complex plane, and represents CP-violation. To understand CP-violation, the angles of this triangle should be measured precisely.

Present limits for each angle

$$\begin{aligned} \phi_1 &= 21.15^\circ \begin{smallmatrix} +0.90^\circ \\ -0.88^\circ \end{smallmatrix} \\ \phi_2 &= 89.0^\circ \begin{smallmatrix} +4.4^\circ \\ -4.2^\circ \end{smallmatrix} \\ \phi_3 &= 68^\circ \begin{smallmatrix} +13^\circ \\ -14^\circ \end{smallmatrix} \end{aligned}$$

The measurement accuracy of ϕ_3 is not so good, and should be improved.

ϕ_3 can be measured by examining the asymmetry between $B^- \rightarrow DK^-$ and $B^+ \rightarrow DK^+$ decays. Among the various B^\pm decays, B^\pm meson which decays to neutral D meson (D^0 or \bar{D}^0) and K^\pm meson is used for ϕ_3 measurement.



D^0 and \bar{D}^0 can decay to the same final states. Therefore $B^- \rightarrow D^0 K^-$ and $B^- \rightarrow \bar{D}^0 K^-$ decay amplitudes interfere each other. The interfering between D^0 and \bar{D}^0 is used to measure ϕ_3 .

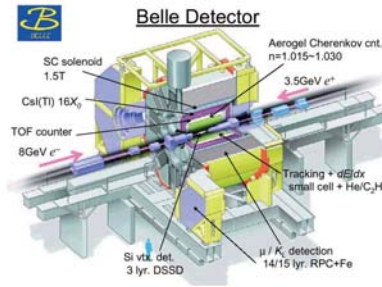
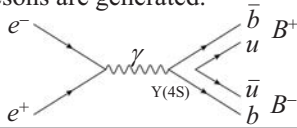
2. Facility

KEKB-factory & Belle Detector

- KEKB-factory is a facility to make B particles.
- High energy electrons and positrons collide, and annihilate in pairs.
- From the pair annihilation, heavy B mesons are generated.



KEK@Tsukuba

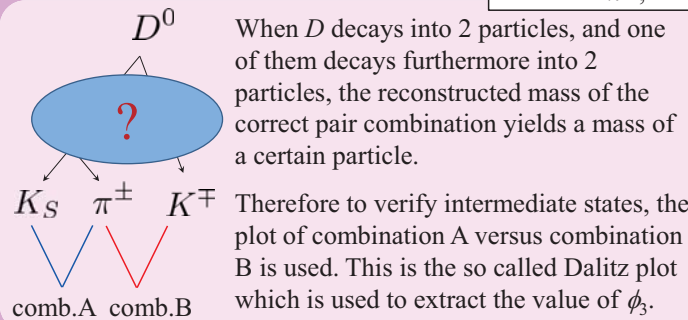


- Belle detector is to search the decays of B particles.
- Belle detector consists of many sub-detectors, and determines the particle type, momentum, charge, and so on.
- The mother particles are reconstructed from detected particles.
- We have recorded the world largest data of 1014 fb^{-1} .

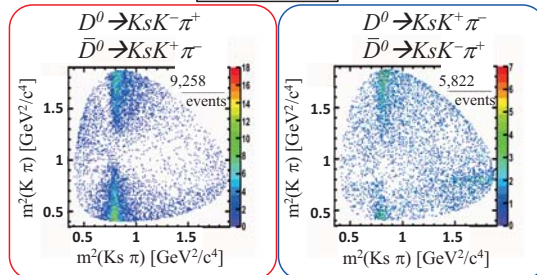
3. Analysis

- Neutral D particles decay to various particles. In this study, $D \rightarrow [K_S K^\pm \pi^\mp]$ decay is searched.
- There are 2 modes in $D \rightarrow K_S K^\pm \pi^\mp$ decays: $D^0 \rightarrow K_S K^- \pi^+$, $D^0 \rightarrow K_S K^+ \pi^-$ and their charge conjugate mode, because both D^0 and \bar{D}^0 can decay into $K_S K^- \pi^+$ and $K_S K^+ \pi^-$.
- D decays into $K_S K \pi$ via many intermediate processes (e.g. $D \rightarrow [K_S \pi^+ K^{*-}] K^-$, $D \rightarrow [K^- \pi^+ K^{*0}] K_S$, ... etc.).
- These processes should be separated because strong phases differ. The Dalitz plot analysis is needed.
- $D^0 \rightarrow K_S K^- \pi^+$ cannot be distinguished from $\bar{D}^0 \rightarrow K_S K^- \pi^+$ in $B^\pm \rightarrow DK^\pm$, however, the information of each Dalitz plane is needed to fit $B^\pm \rightarrow DK^\pm$, $D \rightarrow K_S K^\pm \pi^\mp$ Dalitz plane.
- Therefore, $D^{*\pm} \rightarrow D\pi^\pm$, $D \rightarrow K_S K^\pm \pi^\mp$ decay which has large statistics and can be distinguished between D^0 and \bar{D}^0 using the charge of $D^{*\pm}$ is studied to model the Dalitz distribution of $D \rightarrow K_S K^\pm \pi^\mp$ decay.

$D^{*\pm} \rightarrow D\pi^\pm$, $D \rightarrow K_S K^\pm \pi^\mp$: Dalitz analysis



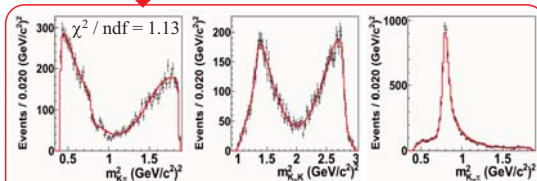
Dalitz Plot



The Dalitz plot of $B \rightarrow DK$ is fitted using the fit result of these Dalitz plots.

2 dimensional fit

141 fb^{-1}



$D^0 \rightarrow K_S K^- \pi^+$ and charge conjugate mode are fitted using 141 fb^{-1} data sample. The another mode is being prepared to fit.

4. Summary and Plan

- The precise measurement for ϕ_3 is important in terms of verification for CP-violation.
- $B \rightarrow DK$ decay is used for the measurement of ϕ_3 .
- Among the various D decays, we use $D \rightarrow K_S K^\pm \pi^\mp$ with Dalitz plot.
- To make the model of $D \rightarrow K_S K^\pm \pi^\mp$ decay, $D^{*\pm} \rightarrow D\pi^\pm$, $D \rightarrow K_S K^\pm \pi^\mp$ is analyzing. $D^0 \rightarrow K_S K^- \pi^+$ have been fitted.
- Of course, the final purpose is the measurement of ϕ_3 using $B^\pm \rightarrow DK^\pm$, $D \rightarrow K_S K^\pm \pi^\mp$.

New geometric interpretation of D-branes and DBI action

Particle Theory and Cosmology Group, Department of Physics

Shuhei Sasa (D1)

In collaboration with Tsuguhiko Asakawa, Satoshi Watamura.

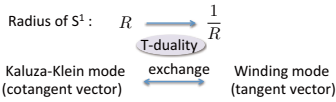
1. Introduction

Symmetries of string theory:

There are various symmetries in string theory. In this poster, we will focus mainly on diffeomorphisms, B-field gauge transformations and T-duality.

T-duality

- Simplest example: S^1 compactification



- General background with U(1) isometry:
for NS-NS flux: a metric g
a Kalb-Ramond 2-form B (called B-field)
a dilaton ϕ

→ T-duality in these background is known as the Buscher rule.

What is a T-duality invariant formulation of SUGRA?

→ Candidate: Generalized Geometry

(There is another candidate; the doubled field theory which is locally same.)

2. Quick review of generalized geometry

A basic object of generalized geometry

Let M be a D -dimensional smooth manifold (for the target space).

The generalized tangent bundle over M ,

$$\mathbb{T}M = TM \oplus T^*M$$

is endowed with

- a fiberwise symmetric non-degenerate bilinear form
 $\langle X + \xi, Y + \eta \rangle = \frac{1}{2}(\iota_X \eta + \iota_Y \xi) = \frac{1}{2} \begin{pmatrix} X^i & \xi_i \\ \eta^j & 0 \end{pmatrix} \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix} \begin{pmatrix} Y^j \\ \eta_j \end{pmatrix}$
- the Dorfman bracket

$$[X + \xi, Y + \eta] = [X, Y] + \mathcal{L}_X \eta - \iota_Y d\xi$$

for $X + \xi, Y + \eta \in \Gamma(TM \oplus T^*M)$.

ι_X : a interior product by X

Symmetries of the symmetric bilinear form $\langle \cdot, \cdot \rangle$

In particular, we will focus on two types of symmetries as follows:

- Diffeomorphisms:**
For a diffeomorphism $f: M \rightarrow M$,
 $f_* \oplus f^*: TM \oplus T^*M \rightarrow TM \oplus T^*M$
is induced. This is called generalized pushforward

- B-transformations:**
For a 2-form $B \in \Omega^2(M)$,
 $e^B: \mathbb{T}M \rightarrow \mathbb{T}M, \quad X + \xi \mapsto X + \xi + \iota_X B$

For an exact 2-form $B = dA \in \Omega^2_{\text{ex}}(M)$, this provides a B-field gauge transformation.

The symmetries of $\langle \cdot, \cdot \rangle$ form a $O(D, D)$ symmetry group.

(There are still other transformations.)

Generalized Geometry

- treats gauge transformations for a Kalb-Ramond 2-form B (called B-field) the same as diffeomorphisms.
- The infinitesimal generators of these transformations are
 - tangent vector fields for diffeomorphisms
 - differential 1-forms for B-field gauge transformations
 → The sections of $\mathbb{T}M \oplus T^*M$ act on itself.
- We define generalized Lie derivatives under these transformations as

$$\mathcal{L}_{X+\xi}(Y+\eta) = \underbrace{\mathcal{L}_X(Y+\eta)}_{\text{Usual Lie derivative for a diffeomorphism}} - \underbrace{\iota_Y d\xi}_{\text{Corresponding to a B-field gauge transformation}}$$

Generalized metric structure (target space)

- We define a generalized metric as a maximal positive-definite subbundle $C_+ \subset \mathbb{T}M$, i.e. for $\forall a, b \in \Gamma(C_+) \setminus \text{zero section}$,

$$\langle a, b \rangle > 0$$

and C_+ has the rank D .

Decompose: $\mathbb{T}M = C_+ \oplus C_-$.
 $C_- = C_+^\perp$ is an orthogonal component of C_+ and has a negative-definite value for $\langle \cdot, \cdot \rangle$.

- The structure group $O(D, D)$ of $\mathbb{T}M$ is reduced to $O(D) \times O(D)$ by specifying the generalized metric C_+ .

- Since the intersection $TM \cap C_+ = \{0\}$, we can describe the generalized metric as a graph of $E = g + B: TM \rightarrow T^*M$:
 $V_+ = v^M(\partial_M + (g + B)_{MN} dx^N) \in \Gamma(C_+)$

- Another definition of the generalized metric is a self-adjoint orthogonal endomorphism $G: \mathbb{T}M \rightarrow \mathbb{T}M$ such that

$$\langle a, Ga \rangle > 0 \quad \text{for } \forall a, b \in \Gamma(C_+) \setminus \text{zero section}$$

$$G^2 = G^T G = 1 \quad \begin{matrix} \longrightarrow \text{Eigenvalues: } G = \pm 1 \\ \longrightarrow C_\pm = \text{Ker}(1 \mp G) \end{matrix}$$

The generalized metric G is given by

$$G = \begin{pmatrix} -Bg^{-1} & g^{-1} \\ g - Bg^{-1}B & g^{-1}B \end{pmatrix}$$

Restriction of the Generalized metric to TM gives a modified Riemannian metric $g - Bg^{-1}B$.

Dilaton structure (target space)

We redefine a dilaton ϕ as a new $O(d, d)$ invariant dilaton d such that

$$e^{-2d} = e^{-2\phi} \sqrt{|\det g|}$$

Strictly speaking, the measure $e^{-2d} dx^0 \wedge \dots \wedge dx^{D-1}$ is $O(d, d)$ invariant.

Then under T-duality,

$$\phi \rightarrow \phi - \frac{1}{2} \ln \det E_{ij}$$

3. D-brane in the framework of generalized geometry

Dp-brane

- is defined as a $(p+1)$ -dimensional hypersurface $S \subset M$ on which open strings can end. Zero modes of an open string living on the D-brane worldvolume S include a gauge field A_a , $(D-p)$ scalar fields Φ^i describing the transverse fluctuations and their superpartners.

- is topologically expressed by (S, φ, V) for the condensed background fields A_a and Φ^i .

$\varphi: S \hookrightarrow M$ is a p -dimensional embedding of submanifold.
 $V \rightarrow M$ is a complex vector bundle with a connection A_a (a so-called Chan-Paton bundle). The rank of V equals the number of coincident D-branes.

(We consider only Abelian gauge theory.)

Questions:

- Dirac-Born-Infeld (DBI) action gives the low-energy effective theory of a D-brane. It can be derived by calculating a partition function in string theory. However it is not obvious why DBI action appears.
- By T-duality, the gauge field is replaced with the scalar field and vice versa. However geometrical meaning of these fields are different. Can we treat these fields alike in the framework of geometry?
- The scalar fields appear as NG boson for broken translational symmetry. How about the gauge field?

Embedding of a Dp-brane worldvolume and a static gauge

Let σ^α : coordinates on a D-brane S ($\alpha = 0, \dots, p$)
 $x^M = (x^a, x^i)$: coordinates on a target space M
($a = 0, \dots, p, \quad i = p+1, \dots, D$)

Physical gauge

Using a diffeomorphism on the D-brane worldvolume S , we choose a static gauge $\sigma^\alpha = x^a$.

The coordinates on S : $(x^a, \Phi^i(x^a))$

Embedding $\varphi: S \hookrightarrow M$ induce the pushforward

$$\varphi_*: TS \rightarrow TM, \quad \varphi_* \left(v^\alpha(\sigma) \partial_\alpha^{(\sigma)} \right) = v^\alpha(x^a, \Phi^i(x^a)) (\partial_a + \partial_a \Phi^i \partial_i)$$

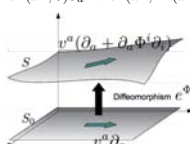
These vector fields can be described by using a diffeomorphism of the target space M such that

$$e^\Phi: S_0 \rightarrow S, \quad (x^a, 0) \mapsto (x^a, \Phi^i(x^a))$$

and its pushforward

$$(e^\Phi)_*: TS_0 \rightarrow TS$$

$$v^\alpha(x^a, 0) \partial_a \mapsto v^\alpha(x^a, \Phi^i(x^a)) (\partial_a + \partial_a \Phi^i \partial_i)$$



A D-brane in the framework of generalized geometry

Dirac structure:

A subbundle $L \subset \mathbb{T}M = TM \oplus T^*M$ over a submanifold $S \subset M$ with the rank D such that

$$\begin{matrix} \text{Isotropic:} & L = L^\perp & \text{i.e. } \langle a, b \rangle = 0 & \text{for } \forall a, b \in \Gamma(L) \\ \text{Integrability:} & [a, b] \in \Gamma(L) \end{matrix}$$

Strategy for constructing DBI action

- We choose a "standard" Dirac structure $L_0 = TS_0 \oplus N^*S_0 \subset \mathbb{T}M$ over S given by sections of the form

$$V_0 = v^a(x^a, 0) \partial_a + \xi_i(x^a, 0) dx^i \in \Gamma(L_0)$$

- We define a Dirac structure $L = e^F(e^\Phi)_* L_0 \subset \mathbb{T}M$ where $F = dA$ is a two-form field strength.

$$V_L = v^a(x^a, \Phi^i(x^a)) (\partial_a + \partial_a \Phi^i \partial_i + F_{ab} dx^b) + \xi_i(x^a, \Phi^i(x^a)) (dx^i - \partial_a \Phi^i dx^a) \in \Gamma(L)$$

- Generalized metric structure

When viewed from the Dirac structure L_0 , since L_0 is isotropic and the intersection $L_0 \cap C_+ = \{0\}$, the generalized metric can also be written as a graph of $t = s + a: L_0 \rightarrow L_0^\perp$,

$$V_L = v^a(\partial_a + t^i_a \partial_i + t_{ab} dx^b) + \xi_i(x^a, t^i_b dx^b + t^{ij} \partial_j) \in \Gamma(C_+)$$

We find the relation

$$\begin{matrix} t_{ab} & = & E_{ab} - E_{ai} E^{ij} E_{jb}, & t^j_a & = & -E_{ak} E^{kj} \\ t^j_b & = & E^{jk} E_{kb}, & t^{ij} & = & E^{ij} \end{matrix}$$

These relations are similar to the Buscher rule since T-duality exchanges vector fields and 1-forms. (But this is not T-duality.)

Restriction of the generalized metric G to L gives a modified metric $s - (a + F)s^{-1}(a + F)$.

$$\left[\begin{matrix} \text{abbreviation: for example, } D \times D \text{ matrix } s = \begin{pmatrix} s_{ab} & s^{ij} \\ s^i_b & s^{ij} \end{pmatrix} \end{matrix} \right]$$

Here,

$$F = \begin{pmatrix} F_{ab} & \partial_a \Phi^i \\ -\partial_b \Phi^i & 0 \end{pmatrix}$$

We have

$$\det^{\frac{1}{2}}(s - (a + F)s^{-1}(a + F))_{MN} = \det^{-\frac{1}{2}} s_{MN} \det^{\frac{1}{2}}(s + a + F)_{MN}$$

- Dilaton structure

We can rewrite the dilaton structure in terms of the symmetric tensor s_{MN} .

$$e^{-d} = e^{-\phi} \det^{\frac{1}{2}} s_{MN} \det^{\frac{1}{2}} E_{ij}$$

- DBI action

Combining the result of 4 and 5, we get

$$\begin{aligned} S_{\text{DBI}} &= \int_S e^{-\phi} \sqrt{|\det E_{ij}|} \sqrt{|\det(s + a + F)_{MN}|} dx^0 \wedge \dots \wedge dx^p \\ &= \int_S e^{-\phi} \sqrt{|\det(\varphi^*(g + B) + F)_{\text{add}}|} dx^0 \wedge \dots \wedge dx^p \end{aligned}$$

4. Non-linear realization

Non-linear realization

- Consider infinitesimal transformation with generators

$$\begin{aligned} \epsilon &= \epsilon_\parallel + \epsilon_\perp = e^a \partial_a + \epsilon^i \partial_i \\ \Lambda &= \Lambda_\parallel + \Lambda_\perp = \Lambda_a dx^a + \Lambda_i dx^i \end{aligned}$$

- By computation of the generalized Lie derivative $\mathcal{L}_{\epsilon+\Lambda} V_+$, we can read that

$$\begin{aligned} \delta t_{ab} &= -e^c \partial_c t_{ab} - \partial_a \epsilon^c t_{cb} - t_{ac} \partial_b \epsilon^c - \partial_a \Lambda_i t^i_b - t^i_a \partial_b \Lambda_j + \partial_a \Lambda_i \\ \delta t^j_a &= -e^c \partial_c t^j_a - \partial_a \epsilon^c t^j_c - \partial_a \Lambda_i t^{ij} + \partial_a \epsilon^j \\ \delta t^i_b &= -e^c \partial_c t^i_b - t^i_c \partial_b \epsilon^c - t^{ij} \partial_b \Lambda_j - \partial_b \epsilon^i \\ \delta t^{ij} &= -e^c \partial_c t^{ij} \end{aligned}$$

$$\delta \Phi^i = \epsilon^i - e^a \partial_a \Phi^i$$

$$\delta A_a = \Lambda_a - e^b \partial_b A_a - \partial_a \epsilon^b A_b - \Phi^i \partial_i A_a$$

We need these terms for non-linear realization of broken symmetries.

Put x, y :

$$\begin{aligned} S_{\text{DBI}} &= \int_S e^{-2\phi} \det^{\frac{1}{2}} \{s - (a + F)s^{-1}(a + F)\}_{MN} dx^0 \wedge \dots \wedge dx^p \\ &= \int_S e^{-2\phi} \det^{\frac{1}{2}-y} s_{MN} \det^x E_{ij} \det^{2y} (s + a + F)_{MN} dx^0 \wedge \dots \wedge dx^p \end{aligned}$$

Diffeomorphism invariance: $\frac{x}{2} + y = \frac{1}{2}$
Vanishing $\det^{\frac{1}{2}-y} s_{MN}$ for non-linear realization: $\frac{x}{2} - y = 0$

We get $x = \frac{1}{2}, \quad y = \frac{1}{4}$

P-15 Electron microscopy and spectroscopy studies of organic molecules inside SWCNT

Md. Mahbubul Haque, Yohei Sato, Masami Terauchi, Tomonari Wakabayashi^A, Kazu Suenaga^B

Date: 2012-Feb-20

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The 4th GCOE International Symposium on "Weaving Science Web beyond Particle-Matter Hierarchy" Feb 20-22, 2012

Introduction

What are polyynes?

1. Polyynes are one-dimensional linear chain of C atoms having alternate single and triple bonds.
2. They are pure sp-hybridized molecules.
3. The terminals of polyynes may contain atoms or group of atoms
4. The simplest form of polyynes contain H atoms at both the terminals.
5. The general formula of simplest polyynes can be written as

General formula



An example of polyyne molecule (C₆H₂)



Basic properties

- Unstable in air or even in liquid at high concentration
- Very stable inside SWCNT

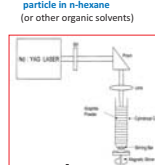
Why polyynes are interesting?

1. Physics: We can observe directly the physical properties of such usually rare one-dimensional material
2. Chemistry: Longer polyynes may be possible to obtain by connecting smaller polyynes hence their chemical properties can be observed.
3. Application: It is expected that polyynes could be the smallest semiconducting molecular nanowire.

Synthesis of polyynes

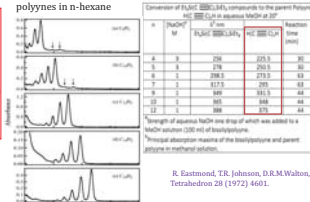
Polyynes can be found naturally as interstellar objects.
In laboratory, synthesis of polyynes can be carried out using several methods.

Laser ablation of graphite particle in n-hexane (or other organic solvents)



HPLC (High Performance Liquid Chromatography) coupled with UV multichannel and VIS spectrometer

UV-vis absorption spectra of polyynes in n-hexane



R. Eastmond, T.R. Johnson, D.R.M. Walton, Tetrahedron 28 (1972) 4661.

H. Tabata et al. Carbon 44, 3168-3176, 2006

Our Motivation

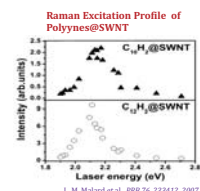
We would like to carry out-

Microscopic studies for polyynes@SWCNT to obtain-
HRTEM images of polyyne to understand their encapsulation, distribution, orientation etc. inside SWCNT

And

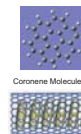
Spectroscopic studies -

EELS (electron energy loss spectra) of polyyne@SWNT-
to understand their electronic properties. We would like to compare EELS data with Raman excitation profile of polyynes trapped inside nanotube

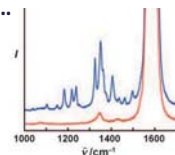


L. M. Malard et al. PRB 76, 233412, 2007

Our Motivation...



Intermolecular spacing is 3.5 ± 0.3 Å



Raman spectra of coronenes@SWCNT (blue)
Raman spectra of SWCNT (red)

We would like to get
Electron diffraction and EELS of coronene@SWCNT in order to study structural and electronic properties of molecular coronenes inside SWCNT.

Experimental Condition

Specimen (C₁₀H₂@SWCNT) is collected from T. Wakabayashi lab at Kinki University

Before the experiments, specimen was heated up until 250 degree for 2 hours in vacuum.

Instrument

TEM equipped with monochromator

HRTEM (Mono OFF)

60kV at room Temperature

EXPtime: 5 sec

Detector: CCD camera

and negative films

EELS (Mono ON)

Accelerating voltage :60kV

at Room Temperature

Probe size: ~1.5 nmφ

Beam current: 2pA

ΔE~65meV

EXPtime: 18-28sec

Detector: Imaging Plate (IP)

Monochromator TEM

Electron source :FEG

Monochromator: Wien-filter

TEM column :JEM-2010FEF

Analyzer :Q-filter

Detector :Imaging Plate

Condition of EELS

HT :60kV

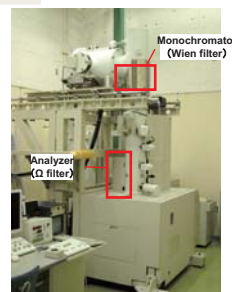
Probe size: 1.5nmφ

Temperature: 100K (LiqN₂ holder)

ΔE :61 meV-0.15 eV

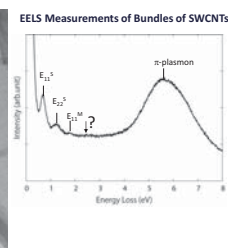
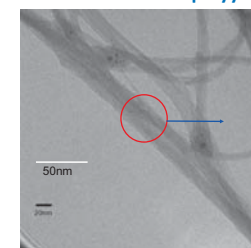
Exp.time: 6s (Low-loss)

256 s (Core-loss)



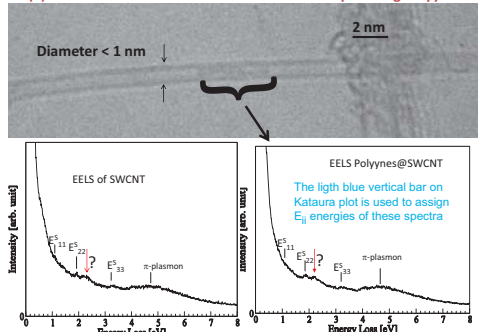
Instrument

Some of our results -of polyynes

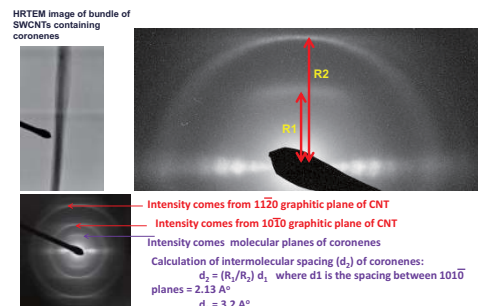


Existence of polyynes inside the bundle of SWCNTs is not clear from this EELS data

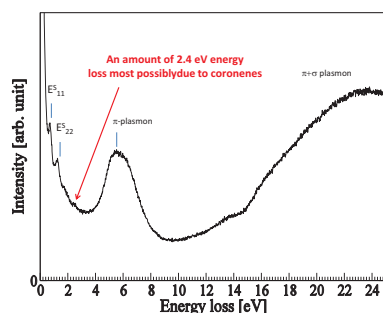
(1) EELS Measurements of Isolated SWCNT encapsulating Polyynes



Some of our results – ED of coronenes@SWCNTs



Some of our results -EELS of coronenes@SWCNTs



Summary and Future Plan

- The encapsulation of polyynes by SWCNT has been confirmed by HRTEM images.
- The distribution of polyynes inside nanotube is not uniform.
- The orientation of polyynes seems bending inside large-diameter nanotube.
- The EELS measurements seems that the electron energy loss due to polyynes depends on diameter of nanotube.

EELS measurement	Diameter of SWCNT in nm (approx.)	Electron loss energy of Polyynes (eV)
1	0.9	2.3
2	1.1	2.55
3	1.2	2.62

- The encapsulation of coronenes by SWCNT is confirmed by electron diffraction from bundle of SWCNT.

Medium heavy Λ hyper nuclear spectroscopic experiment by the $(e,e'K^+)$ reaction



Graduate school of science, Tohoku University
Toshiyuki Gogami for HES-HKS collaboration



1. Introduction



Figure.1 : HES-HKS group photo in the experimental hall C in JLab (2009).

We have been performing Λ hypernuclear spectroscopic experiment by the $(e,e'K^+)$ reaction since 2000 at Thomas Jefferson National Accelerator Facility (JLab). The $(e,e'K^+)$ can achieve 100 keV (FWHM) energy resolution compared to a few MeV (FWHM) by the (K,π) and (π,K) experiments (Table.1). Therefore, more precise Λ hypernuclear structures can be investigated by the $(e,e'K^+)$ experiment. $^7_\Lambda\text{He}$, $^9_\Lambda\text{Li}$, $^{10}_\Lambda\text{Be}$, $^{12}_\Lambda\text{B}$, $^{28}_\Lambda\text{Al}$ and $^{52}_\Lambda\text{V}$ were measured in the experiment at JLab Hall C. In addition, $^9_\Lambda\text{Li}$, $^{12}_\Lambda\text{B}$, $^{16}_\Lambda\text{N}$ were measured in the experiment at JLab Hall A.

Table.1 : The features of $(e,e'K^+)$ reaction comparing to (π,K) and (K,π) reactions.

Reaction	$(e,e'K^+)$	(π,K)	(K,π)
Momentum transfer (Typical)	~ 300 [MeV/c]	~ 300 [MeV/c]	~ 100 [MeV/c]
Λ 's Spin	flip = non-flip Spin dependent structure	non-flip	non-flip
Λ 's from	proton Mirror lambda hypernuclei	neutron	neutron
Beam	High quality, high intensity primary	secondary	secondary
Target	Thin (~ 100 mg/cm 2) (Isotopically enriched)	Thick ($> a$ few [g/cm 2])	Thick ($> a$ few [g/cm 2])
Energy resolution (FWHM)	≤ 500 [keV] Fine structure	$1-3$ [MeV]	$1-3$ [MeV]

2. Experimental Setup (JLab E05-115)



Figure.2 : The experimental setup of JLab E05-115 (2009)

$$\text{Missing Mass : } M_{\text{HY}}^2 = (E_e + M_T - E_{e'} - E_K) - (p_{e'} - p_K - p_{\Lambda})^2 \rightarrow \begin{cases} \Lambda \text{ binding energy} \\ \text{Cross section} \end{cases}$$

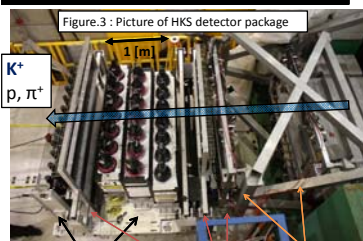
Measure with spectrometers

The $(e,e'K^+)$ experiment is a coincidence experiment between scattered electron and generated kaon. The cross section of $(e,e'K^+)$ are larger at forward scattered angles of these two particles. Therefore, both of these two particles need to be detected at forward angles at the same time. To do so, a dipole magnet (splitter magnet) was set just after the target to separate scattered electron and kaon into different directions.

Figure.1 shows the experimental setup of JLab E05-115. There are HES (High resolution Electron Spectrometer) and HKS (High resolution Kaon Spectrometer) to measure momenta of scattered electron and kaon associated with $(e,e'K^+)$ reaction, respectively. Both HES and HKS consist of QDD magnets.

	HES	HKS
$\Delta p/p$	$\sim 2 \times 10^{-4}$	$\sim 2 \times 10^{-4}$
Momentum [GeV/c]	0.844 ± 0.144	1.20 ± 0.15
Angle [degree]	3.0 – 9.0	1.0 – 13.0
Beam energy [GeV]	2.344	
Target	^7Li , ^9Be , ^{10}B , ^{12}C , ^{52}Cr , $(\text{CH}_2, \text{H}_2\text{O})$ (Hypernuclei)	
	$(^7_\Lambda\text{He}, ^9_\Lambda\text{Li}, ^{10}_\Lambda\text{Be}, ^{12}_\Lambda\text{B}, ^{52}_\Lambda\text{V}) (\Lambda, \Lambda)$	

3. Particle identification



- Cherenkov detectors -AC, WC-
• Aerogel ($n=1.05$)
• Water ($n=1.33$)
- Drift chambers -KDC1, KDC2-
 $\sigma = 250$ [μm]
- TOF walls -2X, 1Y, 1X-
(Plastic scintillators)
TOF $\sigma = 170$ [ps]

There is not only kaon in HKS but also proton and pion as background particles. To measure kaon efficiently, HKS has two drift chambers (KDC1, KDC2) for tracking, three scintillator walls (KT0F1X, 1Y, 2X) for TOF measurement and two type of Cherenkov detectors (WC1, WC2, AC1, AC2, AC3) for proton and pion rejection in the stages of trigger and offline analysis (Fig.3).

Figure.4 shows mass square distribution which can be calculated by the following equation :

$$m^2 = p^2/c^2 - E^2$$

where, p is a momentum reconstructed by a transfer matrix of HKS, and E ($= \sqrt{p^2/c^2 + m^2}$) is derived by TOF measurement. Figure.5 is showing NPE of Cherenkov detectors vs. mass square. We can distinguish pion, kaon and proton with mass square and Cherenkov detector information as you can see in the figure. When the Cherenkov and mass square cut are applied to survive $\sim 90\%$ kaon in the total event, $<2\%$ proton and $<1\%$ pion are contaminated in the kaon.

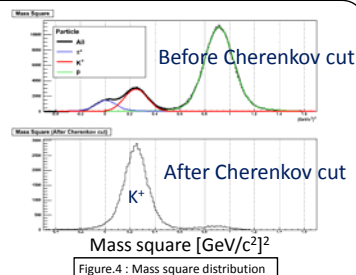


Figure.4 : Mass square distribution

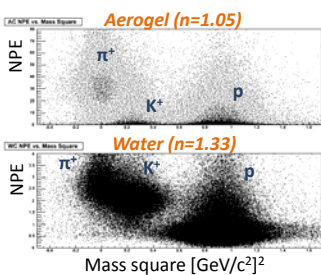


Figure.5 : NPE of Cherenkov detector vs. mass square

4. New tracking code for high multiplicity data

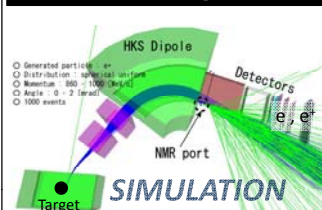


Figure.6 : Background simulation of HKS

We used two planar-type drift chambers which have 6-layers ($uu'xx'vv'$) in each chamber for HKS tracking. Figure.7 shows the multiplicity of the typical layer of the drift chambers. The multiplicity for ^{52}Cr target (~ 5) is much higher than that for CH_2 target (~ 2) as you can see the figure. The conventional tracking code that we used in JLab Hall C cannot handle high multiplicity data efficiently. Therefore, we lose events for high multiplicity data in the tracking stage. To deal with the high multiplicity data, a new tracking code need to be developed.

It is getting harder to perform the $(e,e'K^+)$ experiment as the proton number of the target, Z_{tar} , is larger. This is because that the background events which are caused by electromagnetic processes roughly proportional to Z_{tar}^2 . Mainly, the electron arm suffers from these backgrounds. However, the hadron arm (HKS) also suffers from them which are not on the HKS optics. A positron generated in the target by pair creation process hit the vacuum chamber which is just after the HKS dipole magnet, and generate background events such as positron and electron in the HKS detectors (Fig.6). These background events make the singles rate of HKS be higher ($\sim 30\text{MHz}/\text{plane}$, $8\mu\text{A}$ beam on ^{52}Cr target).

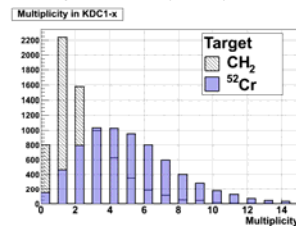


Figure.7 : Multiplicity of the typical layer of HKS drift chambers.

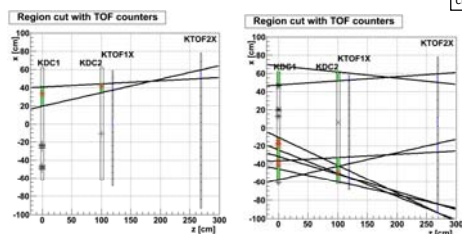


Figure.8 : Event display of the developed hit-wire selection for tracking. Blue squares represent hit TOF segments, green regions are selective regions determined by the TOF combinations, red markers are selected hit-wires for tracking and black markers are hit-wires which are not used for tracking.

A new event selection routine is developed and implemented to increase the tracking efficiency for the high multiplicity data. Before the pattern recognition which is in the first stage of the tracking, hit-wires to use for tracking are selected by the combination of the TOF detectors considering HKS optics (Fig.8).

After this development, the number of kaons is increased by $\sim 130\%$, and also the analysis time is decreased by $\sim 30\%$ for ^{52}Cr target.

5. Energy scale calibration

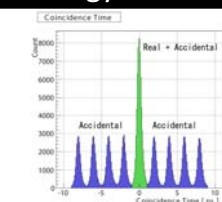


Figure.9 : Coincidence time between HES and HKS.

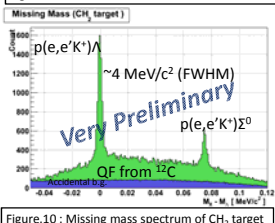


Figure.10 : Missing mass spectrum of CH_2 target

One of the large advantages of the $(e,e'K^+)$ reaction is that the absolute energy calibration can be done with data of Λ and Σ^0 converted from a proton target such as CH_2 and H_2O target. We measure both positions and angles of electron and kaon at the reference planes. Then, those information are converted to momentum vectors at the target by transfer matrices to calculate the missing mass. Therefore, the tuning of the transfer matrices is a heart part to measure hypernuclei with better energy resolution.

Figure.9 shows a coincidence time between HES and HKS after the kaon selection. It is calculated by the following equation :

$$T_{\text{coinc}} = T_{\text{HES}} - T_{\text{HKS}}$$

where, T_{HES} and T_{HKS} are the times at the target which are simply calculated by path length from the focal planes to the target and β of the particles. In the figure, the beam structure of the CEBAF ($\sim 2\text{ns}$ bunch structure) can be seen, and a peak on the center is a bunch which includes real coincidence events.

Figure.10 shows a missing mass spectrum of CH_2 target. When the real coincidence events are chosen, peaks of Λ and Σ^0 clearly can be seen on the Quasi-free Λ events which come from ^{12}C and accidental background events. These peaks are used for the energy scale calibration. Also, thick tungsten alloy sieve slits which are set just before the Q-magnet of each spectrometer are used for calibration of angular component.

6. Summary

- We performed the $(e,e'K^+)$ experiment at JLab Hall-C in 2009 (JLab E05-115), and successfully took data of $^7_\Lambda\text{He}$, $^9_\Lambda\text{Li}$, $^{10}_\Lambda\text{Be}$, $^{12}_\Lambda\text{B}$ and $^{52}_\Lambda\text{V}$.
- Kaon identification
 - When the cut applied to survive $\sim 90\%$, $<2\%$ proton and $<1\%$ pion contaminate in the kaon.
- New tracking code for high multiplicity data (^{52}Cr , H_2O target)
 - The number of kaon is increased by $\sim 130\%$
 - The analysis time is decreased by $\sim 30\%$
- Energy scale calibration is in progress.

Coherent double pion photoproduction on the deuteron

C. Kimura Department of Physics, Tohoku University

Introduction

• Double pion photoproduction

- Investigation of double pion photoproduction from the deuteron is useful in understanding the interaction between photon and bound nucleon.
- Since the deuteron consists of loosely bounded two nuclei, the quasi-free process is considered to dominate the photo-absorption process.
- The impact approximation works well in the meson photoproduction process on the deuteron, that is, the photon interacts mainly with each nucleon in the deuteron (quasi-free process).

• Non quasi-free process

- In the previous study, the cross-section of the double pion photoproduction using NKS(ref.1).
- The reaction through the quasi-free process on the deuteron was found to be smaller than that on the proton. And double delta production was observed and the cross section was obtained (Fig. 1).
- These results suggested the contribution of the non quasi-free process is not small for the photoabsorption on the bound nucleon.

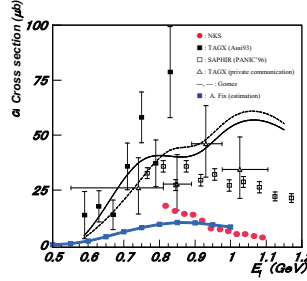


Figure1: Experimental data and theoretical calculation of total cross section of $\gamma d \rightarrow \Delta^{++} \Delta^-$.

• Coherent production

- The process without deuteron breakup is called coherent process(ref.2).
- In the NKS experiment, deuteron was detected in final state. It seems that the contribution of this process to photoabsorption is not small.
- This process was studied by a few theoretical approach and a few experiments (Fig.2).
- We expect to know detailed mechanism of this process using NKS2.

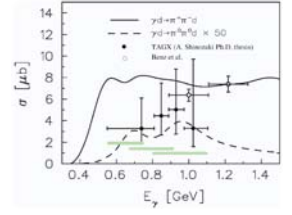


Figure2: Predicted cross section of coherent production and experimental data.

• The aim of this study

- The effect from the binding of the proton and the neutron can be studied in this reaction.
- Investigation of the mechanism of photoabsorption on a bound nucleon using deuteron target and GeV region photon beam.
- Obtaining the cross section and its energy dependence in a wide energy region of photon.
- Investigating the role of coherent photoproduction process ($\gamma d \rightarrow \pi^+ \pi^- d$) in photoabsorption.

Experiment

Photon beam

- The experiment was carried out at Research Center for Electron Photon Science (ELPH), using a photon beam.
- A photon beam is created by bremsstrahlung from an electron beam, and tagged by tagging counters.
- The photon energy range : $0.8 < E_\gamma < 1.1$ GeV (Energy width ~6 MeV)

Neutral Kaon Spectrometer 2 (Fig. 3)

- The charged particles in the final state were detected using the Neutral Kaon Spectrometer 2 (NKS2)
- Dipole magnet (B ~ 0.42 T at the center)
- Drift chambers : Vertex Drift Chamber (VDC), Cylindrical Drift Chamber (CDC)
- TOF counters : Inner Hodoscope (IH), Outer Hodoscope (OH)
- Electron Veto
- Trigger: Detection of more than two charged particles
- Target: Liquid deuterium :located at the center of NKS2



Figure 3 Schematic view of NKS2

Analysis

Event Selection

- Selecting good event using some criteria, hodoscopes timing, tracking χ^2 , vertex distribution, etc.
- Selecting events which include the three charged particle.
- Check the momentum conservation between photon energy and total momentum of three particles.
- Particle identification (Fig.4)
 - Mass of particles
 - Drift Chamber: Reconstructed track
 - Curvature
 - Momentum
 - Hodoscope: Time of flight
 - Velocity of particles

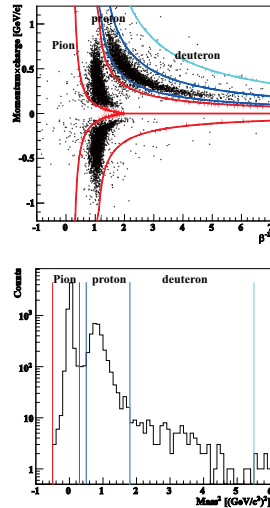


Figure 4: Upper figure is distribution of momentum and inverse velocity of charged particles. Lower figure shows mass square distribution.

- Check the evidence of $\gamma d \rightarrow \pi^+ \pi^- d$ reaction as follows (fig.5):

- Missing mass distribution $\gamma d \rightarrow \pi^+ \pi^- X$.
- Energy deposit at IH.

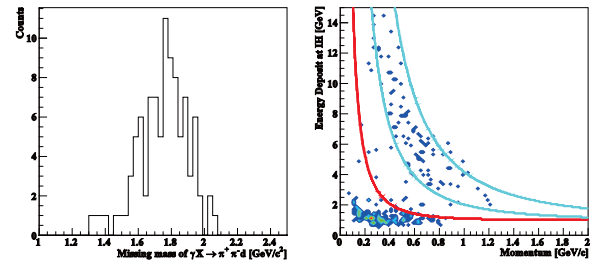


Figure 5: Left is missing mass distribution of $\gamma d \rightarrow \pi^+ \pi^- X$. There is a peak at around deuteron mass (1.88 GeV/c²).

Right shows the distribution of energy deposit at IH. The distribution between two sky blue line which obtained from Bethe-Bloch formula are deuteron. Under the red line are pion.

Summary

- The coherent photoproduction on the deuteron is useful for understanding the mechanism of photoabsorption on a bound nucleon in the GeV region.
- The experiments were carried out in the energy range 0.8-1.1 GeV at ELPH with NKS2.
- The event which include deuteron are detected in our data, and checked these event are $\gamma d \rightarrow \pi^+ \pi^- d$ event.

Further outlook

- Improvement of the analysis method to derive the target reaction.
- Improvement the statics and extend the photon energy region.
- Estimation of the detection efficiencies to obtain the cross section and its energy dependence.
- Discussion the mechanism of this reaction and the effect to photoabsorption.

Reference

- [1] K. Hirose et al.: Phys. Lett. B674 (2009)
- [2] A. Fix and H. Arenhoevel : Eur. Phys.J. A25 115-135(2005)

HKS-HES

The 4th GCOE Internal Symposium

The analysis of Lambda hypernuclear spectroscopic experiment via $(e,e'K^+)$ reaction at JLab Hall-C



Jefferson Lab
Thomas Jefferson National Accelerator Facility

D.Kawama
for JLab E05-115 Collaboration
Department of Physics, Tohoku Univ.



JLab E05-115 experiment : 2009 Aug to Nov at JLab Hall-C

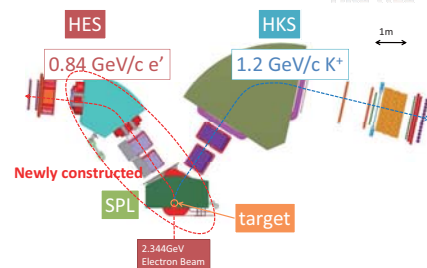
Data Summary

Target	hypernucleus	Target Thickness [mg/cm ²]	Average beam current [μA]	Total Charge [C]
CH ₂	Λ, Σ^0	450.8	2.0	0.28
H ₂ O	Λ, Σ^0	~500	2.7	0.20
⁷ Li	$^7\Lambda\text{He}$	184.0	32.0	4.84
⁹ Be	$^9\Lambda\text{Li}$	188.1	38.3	5.33
¹⁰ B	$^{10}\Lambda\text{Be}$	56.1	38.7	6.25
¹² C	$^{12}\Lambda\text{B}$	112.5	26.8	5.90
¹⁴ Cr	$^{14}\Lambda\text{V}$	134.0	7.6	0.83
		154.0		5.53

- CH₂, H₂O ... for the momentum calibration used as proton target
- $^{12}\Lambda\text{B}$... Study of the structure, linearity check using ground state
 - Charge Symmetry Breaking (CSB) effect, Λ - Σ^0 coupling effect
- $^{10}\Lambda\text{Be}$... Study for medium hypernuclei
 - Fine structure study, Λ single particle energy

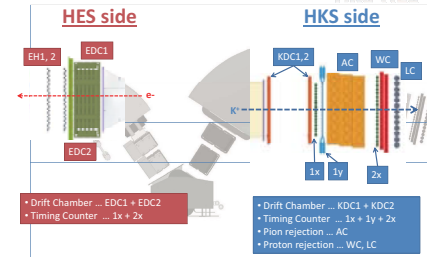
Experimental Setup -magnets-

HKS-HES



Experimental Setup -detectors-

HKS-HES



HES-HKS Spec & Expected energy resolution

High resolution & Wide acceptance spectrometers

	HES	HKS
Configuration	Q-Q-D (50deg)	Q-Q-D (70deg)
e'/K' Central Momentum	0.84 GeV/c	1.2 GeV/c
e'/K' Momentum Acceptance	±0.15 GeV/c	±0.15 GeV/c
e'/K' Solid Angle	~7msr	~8.5msr
Mom. Resolution (FWHM) [MeV/c]	2.0×10^{-4}	2.0×10^{-4}
Ang. Resolution (FWHM) [mrad]	3	2

dE_{beam}=190keV dE_{HKS}=220keV dE_{HES}=190keV

Total Energy Resolution = 350keV (FWHM, design value)

Tracking and kaon PID

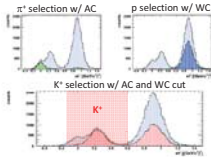
HKS-HES

Tracking ... Performed with enough resolution & efficiencies on both HKS and HES (at least in the low rate data)

Kaon PID

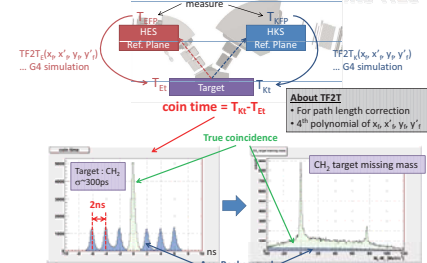
AC ... π^+ , e^+ rejection
WC ... p rejection
 m^2 cut ... clean K^+ selection

Kaon could be selected with ~95% cut efficiency (^{12}C target)



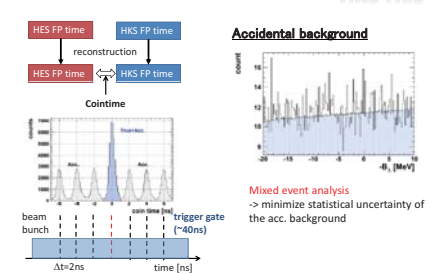
Coin Time Analysis

HKS-HES



Coincidence time

HKS-HES



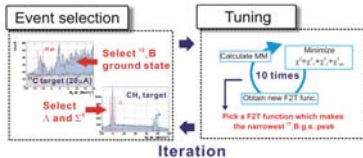
Direct calibration of the F2T function

HKS-HES

$$F2T \text{ function } PK = F2T_{pk}(x_{fk}) \quad 6^{\text{th}} \text{ polynomial of FP quantities}$$

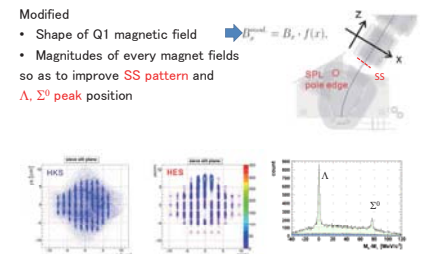
$$= \sum_{i=0}^6 C_i (x_{fk})^i \quad \text{Parameter} \rightarrow \text{Tune}$$

$$\text{Minimizing } \chi^2 \quad \chi^2 = \frac{1}{2} \sum_{i=1}^N \left(\frac{M_i^{\text{data}} - M_i^{\text{F2T}}}{\sigma_i} \right)^2$$



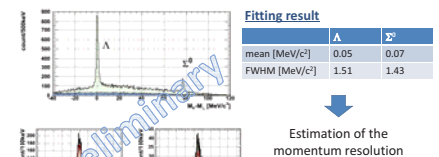
Magnetic field improvement

HKS-HES



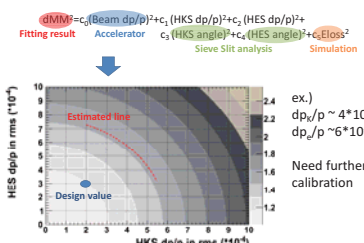
CH2 spectrum

HKS-HES



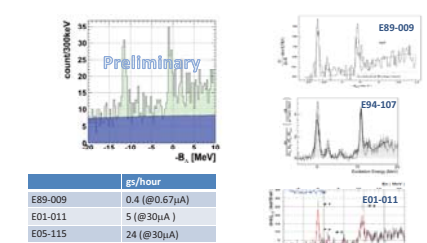
Estimation of momentum resolution

HKS-HES



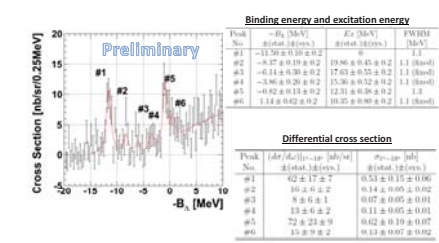
¹²B -spectrum-

HKS-HES



¹²B -fitting results-

HKS-HES



Summary

HKS-HES

JLab E05-115 Experiment, 2009 Aug-Nov
carried out successfully

- Now we are finalizing the data analysis
- Resolution study from the CH₂ spectrum
 - Need further calibration, but the basic way of the calibration was established
- ^{12}B spectrum was consistent with the other data

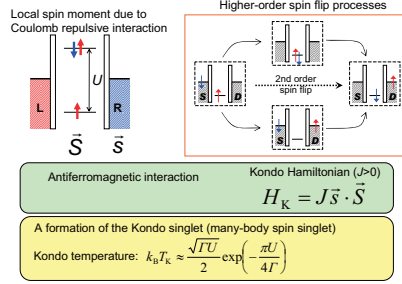
To do

- Tracking analysis for higher rate data
- F2T func. parameter optimization
 - Calibration of the acceptance edge
 - Improvement of the magnetic field

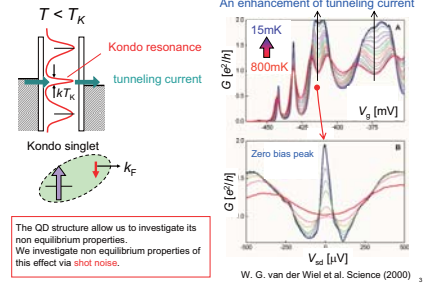
Shot noise measurements for a Kondo-correlated quantum dot in the unitary limit

Yuma Okazaki

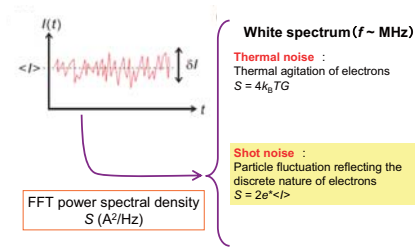
The Kondo effect in a quantum dot



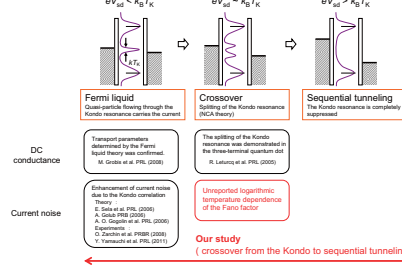
Transport properties of the Kondo quantum dot



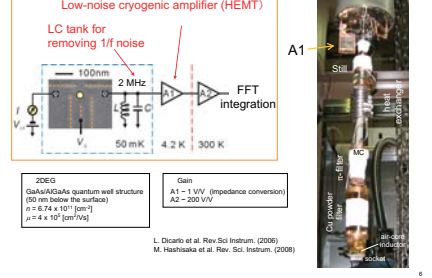
Current noise measurements



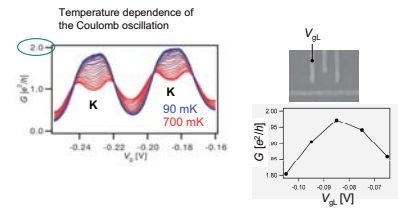
The Kondo effect out of equilibrium



Setup for measuring current noise

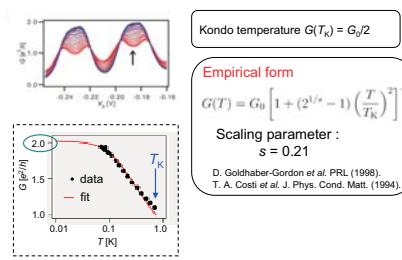


The Kondo effect in the unitary limit

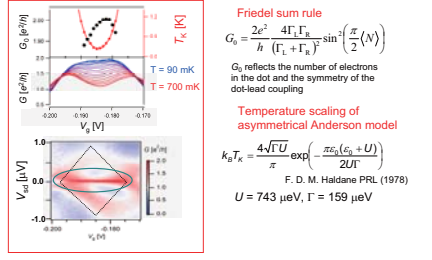


A careful control of the barrier symmetry evolved the Kondo state to the unitary limit, in which uninterested shot noise induced by non-Kondo process can be completely suppressed.

Temperature scaling and Kondo temperature

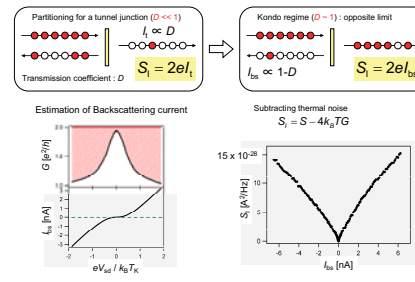


Temperature scaling in the Anderson model

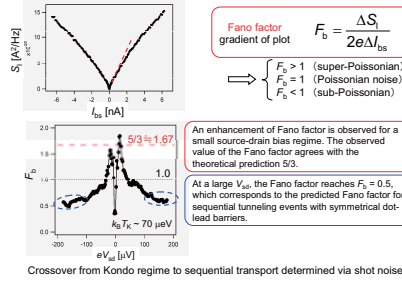


The quantum dot shows characteristic properties of the Anderson model

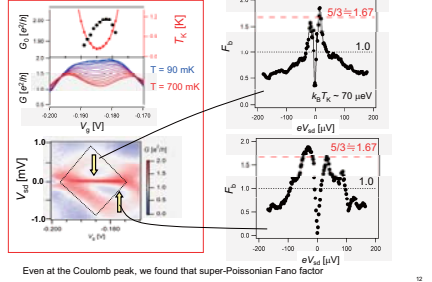
Backscattering current and shot noise



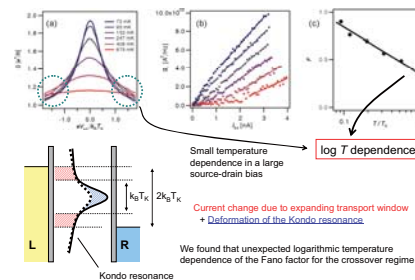
Shot noise measurements and the Fano factor



Out of particle hole symmetry



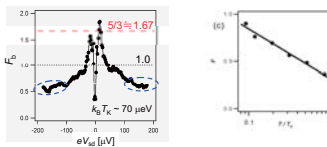
Temperature dependence of the Fano factor



Summary

Shot noise properties of a Kondo correlated quantum dot in the unitary limit

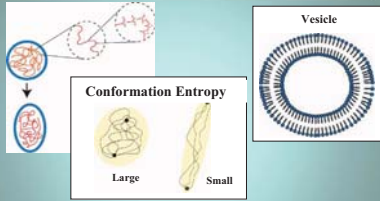
- For a small V_{sd} , the predicted super-Poissonian Fano factor was observed.
- A crossover from Kondo regime to sequential tunneling regime determined via shot noise.
- Super-Poissonian Fano factor was observed even near a Coulomb peak.
- The temperature dependence of the Fano factor in the crossover regime shows unexpected logarithmic temperature dependence.



Equilibrium States of Polymer-Containing Micells

Department of Physics
Yutaka Oya Toshiihiro Kawakatsu

Anisotropic Deformation of Vesicle Induced Polymers



Aim of Our Study
Representation of Vesicle Containing Polymers
Serch for Origin of Anisotropic Deformation

Phase Field Theory

Interface Thickness : ϵ

Helfrich bending energy

$$F_{bend} = \frac{\kappa}{2} \int \left(\frac{1}{R_1} + \frac{1}{R_2} \right)^2 da \rightarrow \frac{\kappa}{2} \int \frac{1}{\epsilon^3} (-\psi + \psi^3 - \epsilon^2 \nabla^2 \psi)^2 d\mathbf{r}$$

R_1, R_2 principle curvature

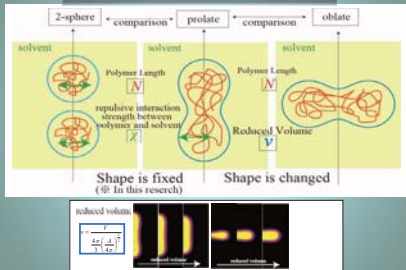
$$F_v = \gamma(V_p - V_0)$$

$$F_A = \sigma(A - A_0)$$

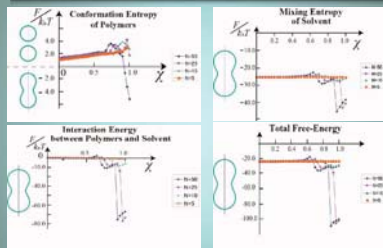
$$\psi = -\frac{\delta(F_{bend} + F_v + F_A)}{\delta\psi}$$

minimize F_{bend}

Parameter of Calculation



Stability between Prorate and 2-Sphere II



When chain length become long, prolate shape become more stable

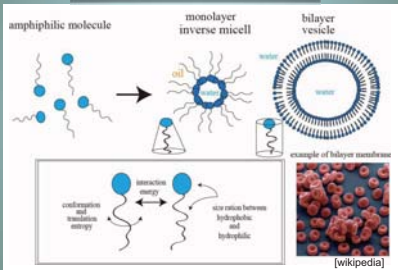
CONCLUSION

Representation of Coupled System of Polymer and Membrane
with Phase Field Theory and Self Consistent Field Theory.

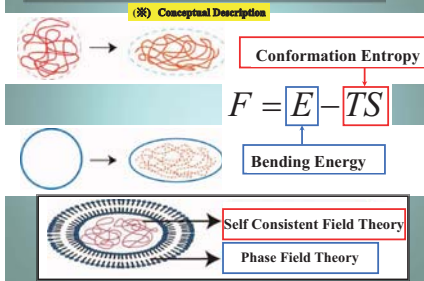
Prolate Shape is More Stable than Oblate Shape,
when Confined Strength of Polymer is Stronger.

Prolate Shape can be More Stable than 2-Sphere Shape,
when Replasive Strength between Polymer and Solvent is exist.

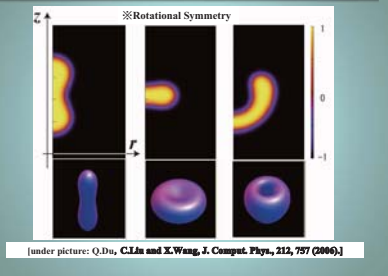
Amphiphilic Membrane



Equilibrium Shape of Vesicle Containing Polymers

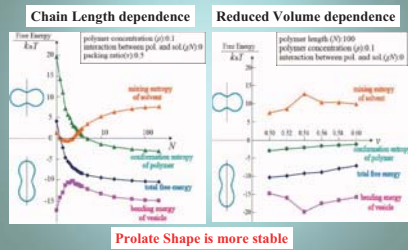


Shape of Amphiphilic Membrane with PF model

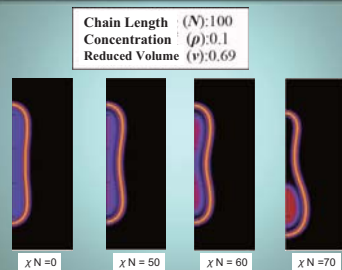


[under picture: Q.Du, C.Lin and X.Wang, J. Comput. Phys., 212, 797 (2006).]

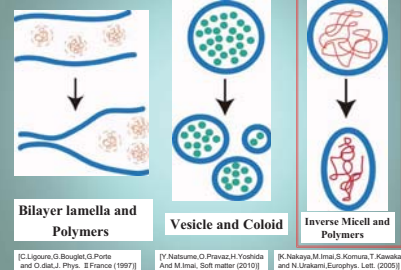
Stability between Prolate and Oblate



Shape Change of Prolate Membrane Induced by Polymer Distribution



Anphiphilic Membrane and Guest Particles

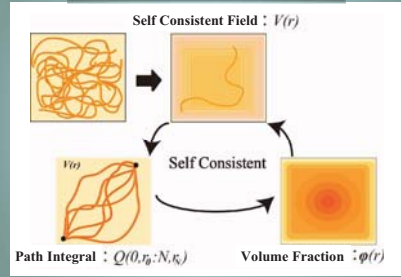


[C. Lipman, G. Boullet, G. Porta and O. D. J. Phys., 18 France (1997)]

[Y. Nakamura, O. P. P. H. Yoshida and M. Imai, Soft matter (2010)]

[K. Nakaya, M. Imai, S. Komura, T. Kawakatsu and N. Uchikawa, Europhys. Lett., (2005)]

Self Consistent Field Theory



Coupled System of PFT and SCFT

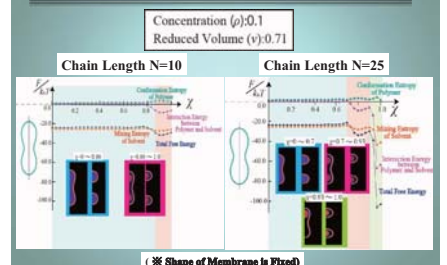
$F(\psi, \phi) = F_{PF}(\psi) + F_{SCF}(\phi) + F_{coupl}(\psi, \phi)$

$$\frac{\kappa}{2} \int \frac{1}{\epsilon^3} (-\psi + \psi^3 - \epsilon^2 \nabla^2 \psi)^2 d\mathbf{r} + \gamma(V_p - V_0) + \sigma(A - A_0) + \int \mu(\mathbf{r}) (\phi_p + \phi_s - \phi_0 - 1) d\mathbf{r}$$

Local Area : $a(\mathbf{r}) = \frac{3\sqrt{3}}{4} \left(\frac{\epsilon}{2} |\nabla \psi|^2 + \frac{1}{4} (1 - \psi^2)^2 \right)$

const. $\times a(\psi)$

Stability between Prorate and 2-Sphere I



Energy Loss of Confinement

Prolate : R_p Oblate : R_o Cylinder (Surface area and Volume are Const) : $R_p \approx 4R_o$

Energy loss of Polymer in Slit

$$F_{Loss} = k_B T \log \frac{N}{2^{1/2}}$$

$$F_{Loss}^{Pr} - F_{Loss}^{Ob} = M \times \left(2 k_B T \log \frac{N}{2^{1/2}} - 1 k_B T \log \frac{N}{2^{1/2}} \right) \propto \frac{V\rho}{R_o^2}$$

(Number of Polymers : $M = \frac{V\rho}{N_{Ad}}$)

Prolate Shape is Stable

“Properties of proton-rich unstable nuclei and two-proton radioactivity”

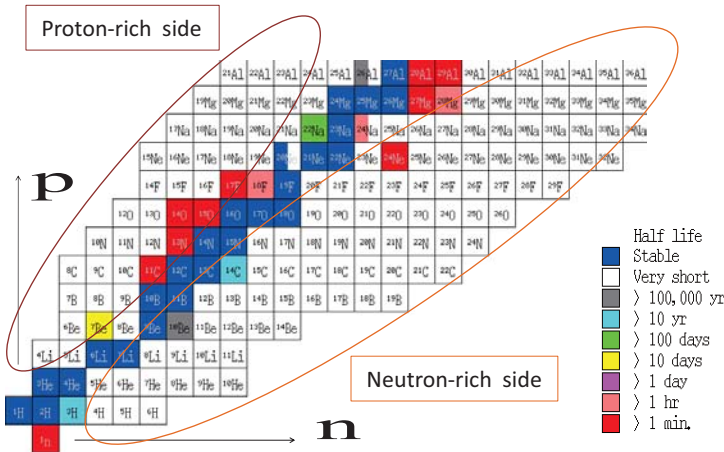
Tomohiro Oishi^A, Kouichi Hagino^A, Takahito Maruyama^A, Hiroyuki Sagawa^B

^A Department of Physics, Tohoku University, ^B Center for Mathematical Sciences, University of Aizu

I. Unstable Nuclei

Nuclei that (1) decay with a quite short lifetime, (2) have a large proton- or neutron-excessiveness, are called “**unstable nuclei**”.

They have some exotic physical properties which are not present in stable nuclei.

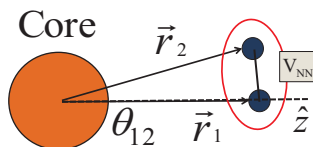
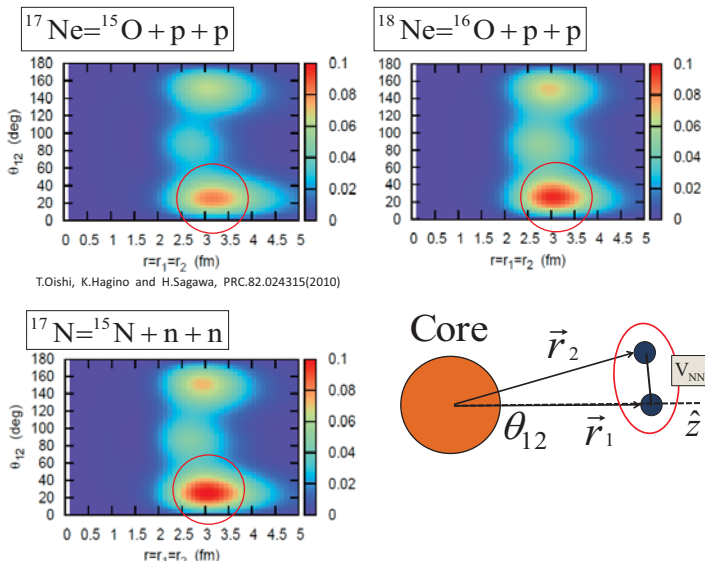


II. Effect of Coulomb Force (EoC)

For the **proton-rich** unstable nuclei, **Coulomb force** (repulsive) is also indispensable besides the nuclear force (attractive).

II(a). EoC for “di-nucleon correlation”

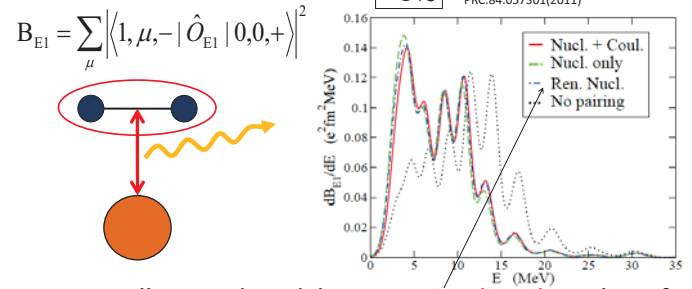
di-nucleon correlation; the spatial localization of two neutrons or protons in nuclei.



	¹⁸ Ne	¹⁷ Ne	¹⁷ N
$\langle V_{NN, \text{Nucl.}} \rangle$ (MeV)	-4.22	-4.01	-3.60
$\langle V_{NN, \text{Coul.}} \rangle$ (MeV)	0.53	0.51	-
S_{2N} (MeV)	-4.57	(-0.94)	9.05
$S_{2N, \text{exp}}$ (MeV)	-4.52	-0.94	8.37

II(b). EoC for “soft dipole excitation”

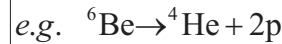
soft dipole excitation; the excitation at low energy of weakly bound nuclei.



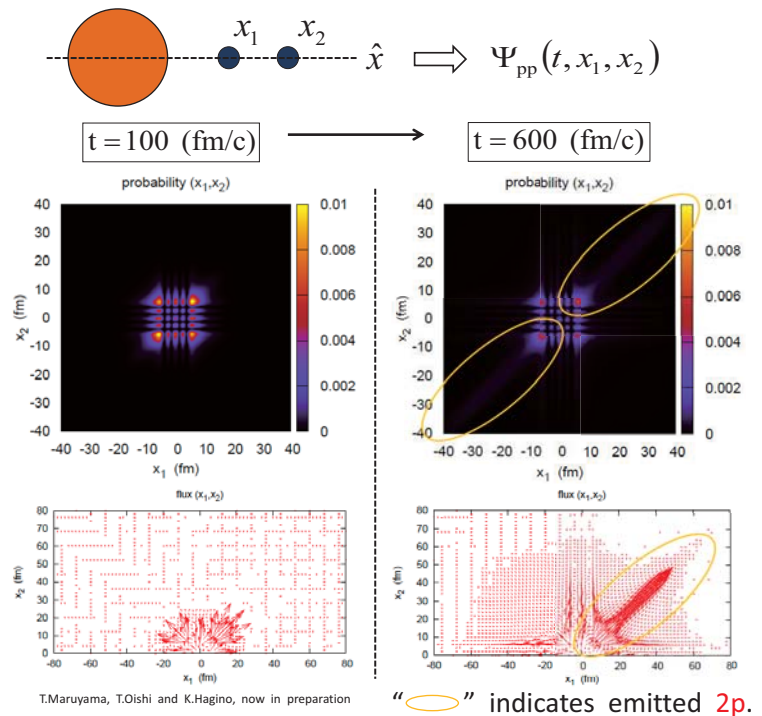
EoC is well reproduced by a **10%-reduced** nuclear force.

III. Two-proton Radioactivity

One of the most important physics of proton-rich nuclei is “**two-proton- (2p-) decay or radioactivity**”.



Analysis with a schematic 1D-3body-model



“**○**” indicates emitted 2p.

IV. Summary

- 1) EoC to the di-nucleon correlation or the soft dipole excitation is investigated. → Coulomb repulsion reduces the pairing interaction by about 10%.
- 2) The time-dependent method is applied to the 2p-decay in a schematic 1D-model. → Time evolution of many-particle-system is well described.

Future work;
realistic 3D-time-dependent analysis of 2p-radioactivity.

Oscillatory Instability of Slow Crack Propagation in Rubbers under Large Deformation

Daiki Endo

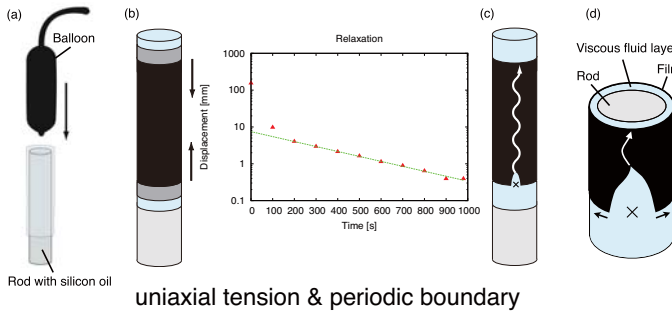
Department of Physics, Tohoku University

Introduction

When a rubber balloon is ruptured, oscillatory crack patterns on a macroscopic scale are observed quite consistently. This oscillating mechanism is still not understood due to the difficulty in measuring rapid crack propagation and treating rapid crack dynamics with significant inertia.

Hence, if we could perform an experiment in which the crack propagation speed is significantly reduced, it may be possible to develop an alternative approach that overcomes those problems.

Experiments of Slow Fracture



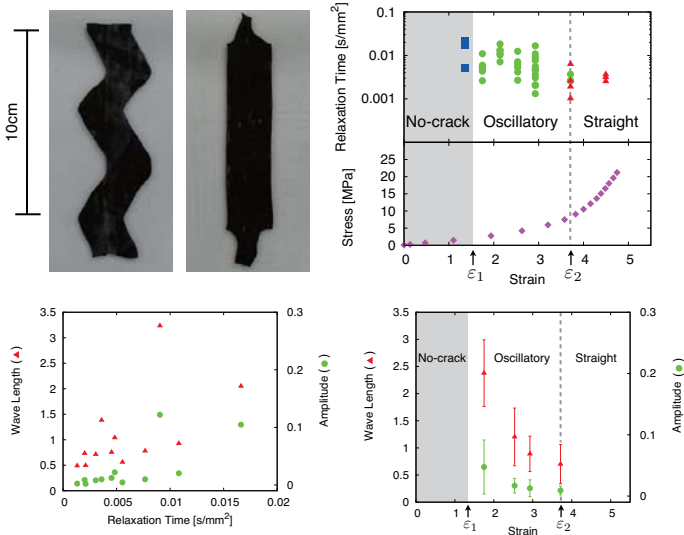
Sample: Qualatex 260Q made by Pioneer Balloon Company		Kinematic viscosity of silicon oils 1,000 ~ 10,000cSt	
Material	Natural rubber	Diameter of rods 15 ~ 35mm	
Thickness	0.31mm		
Circumferential length	20mm		

The above figure shows schematic illustrations of experimental procedure. (a) Oil-coated cylinder is covered by an inflated balloon. By this process, we prepare the rubber film which is stretched on the viscous fluid layer around the cylinder. (b) Relaxation to pure uniaxial tension. (c) Crack initiation. (d) Cross-section of experimental system. The high-viscosity oil between the film and the cylinder significantly reduces the sliding speed of the film with respect to the substrate.

crack speed v 1 mm/s ~ 1 m/s
strain ε 1 ~ 5

$$\rho \ddot{u}_i = \frac{\partial \sigma_{ik}}{\partial x_k} \rightarrow \gamma \dot{u}_i = \frac{\partial \sigma_{ik}}{\partial x_k}$$

- Oscillatory and straight patterns are observed.
- The oscillatory instability is observed at strong nonlinear elastic region.
- Shapes of oscillatory patterns are dependent on the applied strain and the relaxation time.



Summary

- Oscillatory and straight crack patterns are observed in our slow fracture experiments of rubbers.
- The oscillatory instability is observed at strong nonlinear elastic region of rubbers.
- Oscillatory and straight crack patterns are obtained in neo-Hookean model.
- In linear elastic model, oscillatory patterns are not obtained.



The nonlinear elasticity plays an important role in the oscillatory crack pattern formation.

Numerical Simulations

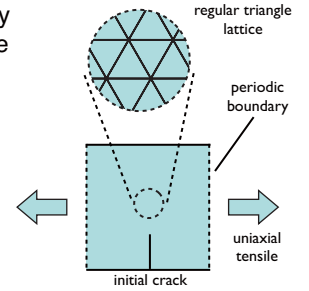
The experimental system is approximated by 2-D system consists of triangle continuum elements

which have { nonlinear elasticity
viscous resistance

with { uniaxial tensile
periodic
initial crack

boundary condition

with appropriate fracture criterion.



energy density function

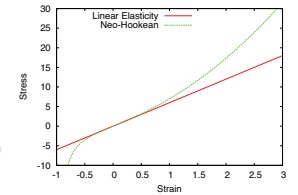
$$F = S \cdot f \begin{cases} \text{Linear Elasticity (Reference)} \\ f = \frac{1}{2} \lambda \varepsilon_{ii} \varepsilon_{jj} + \mu \varepsilon_{ij} \varepsilon_{ij} \quad \lambda = 4, \mu = 2 \\ \text{Neo-Hookean} \\ f = \frac{1}{2} \mu (\lambda_1^2 + \lambda_2^2 + \lambda_3^2 - 3) \quad \mu = 2, \lambda_1 \lambda_2 \lambda_3 = 1 \end{cases}$$

dissipation function of viscous resistance

$$G = \frac{1}{2} \eta \int \left(\frac{dx}{dt} \right)^2 dS$$

equation of motion

$$\frac{\partial}{\partial q_j} \left(\sum_i F_i \right) + \frac{\partial}{\partial \dot{q}_j} \left(\sum_i G_i \right) = 0$$

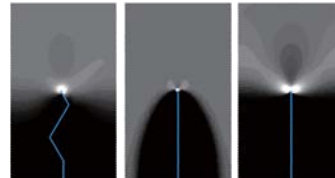


If a strain energy density of i -th element exceeds critical value f_c , the element is removed irreversibly

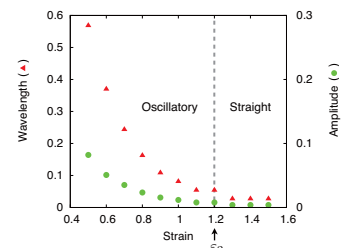
- Oscillatory and straight patterns are obtained in neo-Hookean model. Shapes of the oscillation patterns have a similar strain dependency to the experimental result.
- In linear elastic model, oscillatory patterns are not obtained.

$$f_c = 10 \quad N = 128 \times 256$$

(a) $\varepsilon = 0.5$ (b) $\varepsilon = 1.5$ (c) $\varepsilon = 0.5$



Neo-Hookean Linear



Structure of Neutron-rich Nucleus ^{31}Ne Deduced from Nuclear Reactions

Yasuko Urata, Kouichi Hagino

Department of Physics, Tohoku University

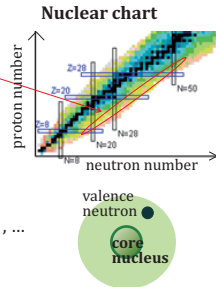
① Halo nucleus

Neutron-rich nuclei

- are not stable against β decays and do not exist in nature.
- recent experimental techniques enable to synthesize and study such nuclei
- have interesting features, such as halo structure

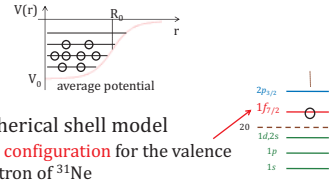
Halo nuclei

Spatially extended density distribution of the weakly bound valence neutron(s): ^{11}Li , ^{11}Be , ^6He , ^{19}C , ...
Halo structure is constructed by a **s- or p-orbit**, for which the centrifugal barrier is small.



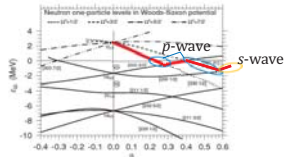
Mean Field Picture

Nucleons are independent particles which are moving in some average potential created by all the nucleons in the nucleus.



Naive spherical shell model
→ $1f_{7/2}$ configuration for the valence neutron of ^{31}Ne

Single particle motion in a deformed potential

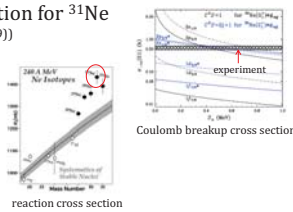


(I.Hamamoto, Phys.Rev.C 81, 021304(R) (2010))

Nilsson diagram
→ $\beta \sim 0.2-0.3$: $[330\ 1/2]$
→ $\beta \sim 0.4-0.6$: $[321\ 3/2]$ } p-wave halo
($\beta \sim 0.6$: $[200\ 1/2]$ → s-wave halo)

Recent Experiment for ^{31}Ne

- Large Coulomb breakup cross section for ^{31}Ne (T.Nakamura, et al., Phys. Rev. Lett. 103, 262501 (2009))
→ soft E1 (electric dipole) excitation
- Large reaction cross section for ^{31}Ne (M.Takechi et al., Nucl.Phys.A 834, 412c (2010))
→ large radius
- ^{31}Ne is a halo nucleus



② Motivation

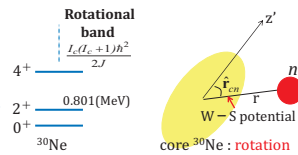
In the Nilsson model, the rotational excitation energy of the core nucleus is neglected. In reality, however, ^{30}Ne has the first 2^+ excited state at 0.801 MeV.

- Inclusion of the excitation energy of the core nucleus with a particle-rotor model
- Calculation of Coulomb breakup cross sections and reaction cross sections

- The effect of finite excitation energy
- Ground state properties of ^{31}Ne

③ Particle-rotor Model

- ^{31}Ne : (deformed core ^{30}Ne) + n
- Axial quadrupole deformation of the core
- The core nucleus has a rotational band.
→ $H = H_{\text{spherical}} + H_{\text{deformed}} + H_{\text{rotation}}$

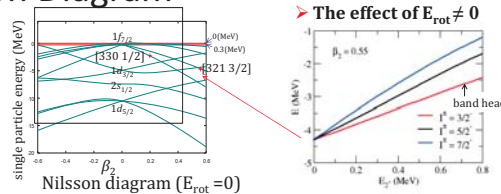


$$\Psi_{IM} = \sum_{I_c J_c} \alpha_{I_c J_c} \left(\begin{matrix} 31\text{Ne} \\ I_c J_c \end{matrix} \right)_{IM} = \sum_{I_c J_c} \alpha_{I_c J_c} \left(\begin{matrix} 30\text{Ne} \\ I_c J_c \end{matrix} \right)_{IM} \otimes \left(\begin{matrix} n \\ 1/2 \end{matrix} \right)_{IM}$$

For example, $0^+ + 1/2^- + 2^+ + 3/2^- + 4^+ + 5/2^- + \dots$

PRM and Nilsson Diagram

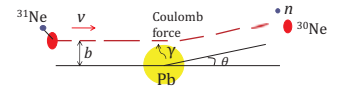
$H_{\text{rotation}} = 0$
→ Nilsson diagram (ϵ_k)



Coulomb breakup

The projectile ^{31}Ne is excited by the electromagnetic field made by the Pb target.

- Breakup into $^{30}\text{Ne} + n$
- Dominant E1 excitation



$$B(E1; i \rightarrow f) = \frac{1}{2I_i + 1} \langle \mathcal{D} \| \hat{D} \| i \rangle^2 \quad \text{Dipole operator is } \hat{D}_\mu = e Z_{\text{eff}} r Y_{1\mu}(\hat{r}_n), Z_{\text{eff}} = -\frac{Z}{A+1}$$

$$\text{Total cross section is } \sigma(E) = \sum_i \frac{16\pi^3}{9hc} N_{E1}(E_f - E_i) B(E1; i \rightarrow f) \quad N_{E1}: \text{virtual photon number}$$

Reaction Cross Section

Reaction cross section (σ_r) is the total cross section except for the elastic scattering.

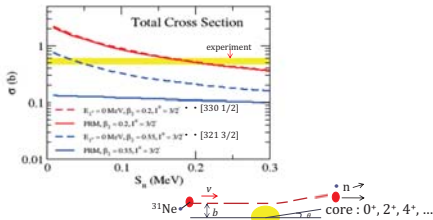
$$\sigma_r = \int d\mathbf{b} \frac{1}{2I+1} \sum_M \langle \Psi_{IM} | (1 - |S|^2) | \Psi_{IM} \rangle \quad |\Psi_{IM}\rangle: \text{initial state wave function (ground state)}$$

S : scattering matrix → $|S|^2$ is the probability that the projectile remains in the elastic channel after the collision with the target.

④ Result: Coulomb Breakup

Y. Urata, K. Hagino, and H. Sagawa Phys. Rev. C 83, 041303(R) (2011)

Total Cross Section

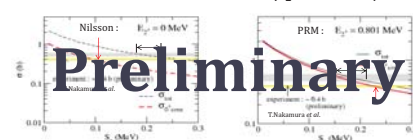


The effect of rotational excitation of core nucleus

$I^\pi = 3/2^-$ configuration at $\beta_2 = 0.55$ cannot reproduce the experimental cross section.

→ $I^\pi = 3/2^-$ configuration at $\beta_2 = 0.2$ is a candidate for the ground state of ^{31}Ne

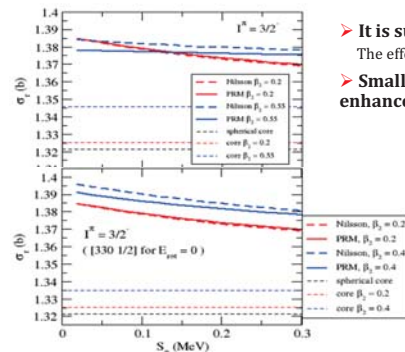
Cross section with 0^+ core • $\beta_2 = 0.2$, $I^\pi = 3/2^-$



Nilsson model calculation
→ does not reproduce the two kinds of data simultaneously.

PRM calculation
→ reproduces both the data consistently.

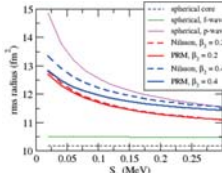
⑤ Result: Reaction Cross Section



Component ($S_n = 0.2\text{ MeV}$)

	$0^+ \otimes p_{3/2}$	$2^+ \otimes p_{3/2}$	Total $p_{3/2}$
$\beta_2 = 0.2$			
Nilsson	24.1%	24.1%	48.2%
PRM	48.5%	8.4%	
$\beta_2 = 0.4$			
Nilsson	26.9%	26.9%	53.8%
PRM	45.6%	13.6%	

Rms Radii



- The behaviors of reaction cross sections are similar to that of rms radii,
- and can be understood from the amount of p-wave component in their configurations.

⑥ Summary

- Ground state properties of ^{31}Ne → particle-rotor model
- rotational excitation energy of the core nucleus

Coulomb breakup cross section

- The effect of excitation energy of the core nucleus
→ $I^\pi = 3/2^-$ configuration at $\beta_2 = 0.55$ is excluded.
- Total and 0^+ cross sections of the Coulomb breakup
→ $I^\pi = 3/2^-$ configuration at $\beta_2 = 0.2$ reproduces the data consistently.

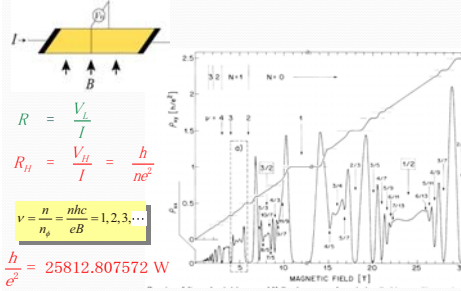
Reaction cross section

- Nilsson model works reasonably well.
- Reaction cross sections and rms radii
→ can be interpreted from neutron configurations.
- Study of β_2 dependence

Investigating of the topological order in the quantum Hall effect

Toru Ito
Department of Physics, Tohoku University

Fractional quantum Hall effect



Laughlin's wavefunction

$$\psi_{Lf} = \prod_{i < j} (z_i - z_j)^m e^{\sum_i \frac{-|z_i|^2}{4\ell^2}}$$

$z = x - iy$: Complex coordinate

$\ell = \sqrt{\hbar / eB}$: Magnetic length

m : Odd integer determined by filling fraction

Numerically, Laughlin's state have over 99% overlap (very good model of the ground state wavefunction)

$\vec{r}_{ij} = z_i - z_j$: Polar coordinate

$$\prod_{i < j} (z_i - z_j)^m = \prod_{i < j} |r_{ij}|^m e^{im\theta_{ij}}$$

This factor represent m magnetic fluxes

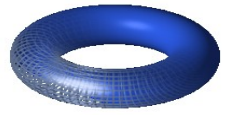
AB effect of quantized magnetic flux



$$\frac{e}{\hbar} \oint \vec{A} d\vec{s} = 2\pi \longleftrightarrow e^{i(m-1)(2\pi)}$$

$$\vec{B} = \begin{pmatrix} 0 \\ 0 \\ \delta(\vec{r}) \frac{\hbar}{e} \end{pmatrix} \quad \vec{A} = \begin{pmatrix} A_r \\ A_\theta \\ A_z \end{pmatrix} = \begin{pmatrix} 0 \\ \frac{\hbar}{e|r|} \\ 0 \end{pmatrix}$$

$\nu=1/m$ Laughlin's state in 1-hole torus geometry has m -fold degeneracy



In g-hole geometry (g-genus Riemann surface), Laughlin's state has m^g -fold degeneracy

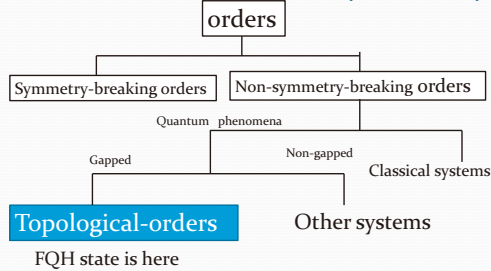
This degeneracy called topological degeneracy

Kinds of Quantum Numbers

- Quantum numbers related to **symmetry**
 - Example: angular momentum \leftrightarrow rotational symmetry
 - Degeneracy from the algebra of the group generators
- Quantum numbers determined by **topology**
 - Examples:
 - Quantized circulation in superfluid He
 - Magnetic flux quantization in superconductors
 - Related to the winding number of a condensate wave function
 - Survive for adding relatively strong perturbation

Hall conductance of FQH is quantum number of this kind

Classification of states (X.G.Wen)



How confirm realization of topological orders?

A new approach is proposed.
This approach use the entanglement entropy of the two divided systems

Definition of the entanglement entropy (EE) of two sub systems

Density matrix of total system : $\rho_{tot} = |\psi\rangle\langle\psi|$

Density matrix of Sub system A : $\rho_A = \text{tr}_B \rho_{tot}$

EE of sub system A : $S_A = -\text{tr}_A \rho_A \ln \rho_A$

EE represents quantum entanglement of two systems.

Example of EE: spin-1/2 system composed of two particles

$$(i) |\Psi\rangle = \frac{1}{2} [|\uparrow\rangle_A + |\downarrow\rangle_A] \otimes [|\uparrow\rangle_B + |\downarrow\rangle_B]$$



Not Entangled

$$S_A = 0$$

$$(ii) |\Psi\rangle = [|\uparrow\rangle_A |\downarrow\rangle_B + |\downarrow\rangle_A |\uparrow\rangle_B] / \sqrt{2}$$

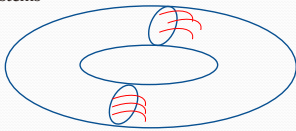


Entangled

$$S_A = \log 2$$

The scaling law of entanglement entropy

Gapped systems have finite correlation length, then the EE is proportional to the boundary length L of two divided systems



$$S_A \propto \alpha L + O\left(\frac{1}{L}\right)$$

In topologically ordered systems, the scaling law of the EE has universal correction γ

$$S_A = \alpha L - \gamma + O\left(\frac{1}{L}\right)$$

γ is called topological entanglement entropy (TEE)

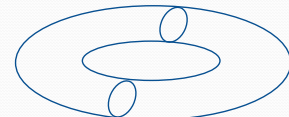
Topological entanglement entropy reflect topological character of the ground state

$$m=3 \text{ Laughlin state: } -\gamma = -n \ln \sqrt{3} = -0.54n$$

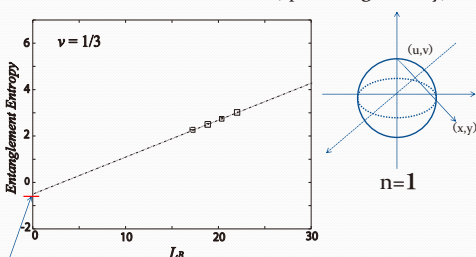
$$m=5 \text{ Laughlin state: } -\gamma = -n \ln \sqrt{5} = -0.805n$$

n =(number of boundaries)

1-hole torus : $n=2$



Numerical calculation of EE(spherical geometry)



Numerically obtained TEE is very close to the theoretical prediction

Summary

- Topologically system is characterized by non-symmetry breaking order
- Such order called topological order
- One of a way of numerically confinement of the topological order is calculating of the topological entanglement entropy

Autonomous decision-making and informed consent

Haruka Hikasa

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1. Introduction

In a time of technical progress, medical technology has made particular progress. Decision making is required in the use of technology. In years gone by, doctors would make the choices regarding the best medical treatment options for the patient. Now, generally the patients themselves make the choice of whether to undergo certain medical treatments, and may even refuse medical treatment at the cost of their life. This is based on a principle of "respect for autonomy"; the acceptance that a person should make decisions according to his/her own values and beliefs. In order to protect that principle, the informed consent of the patient for his/her medical treatment is thought to be indispensable.

However, the relation between autonomy and informed consent is not immediately clear, and various arguments exist concerning it. Therefore, I would like to consider the prerequisites of "respect for autonomy" in decision-making, and the way in which informed consent ought to be practiced. This problem is important to a considered understanding of human life in an age of science and technology.

2. Principle of "Respect for autonomy" in bioethics

In bioethics, the respect for autonomy is one of the most important principles. According to this principle, the healthcare providers have the duty to respect a patient's autonomy, and the providers has the duty to support autonomous choice so that a patient can attain to their own purpose. In other words, healthcare providers disclose relevant information to a patient, and a patient makes a decision having understood the situation. And generally it is thought that informed consent is required in order to apply a principle of respecting patient's autonomy at the place of practice.

3. Compromising a person's autonomy

According to Taylor, however, respect for autonomy cannot serve as the foundation for securing informed consent. "Compromise of a person's autonomy" is diminution of autonomy by being forced by others in decision-making regarding action and choice. Moreover, the autonomy declines with a person being deceived, or acting without information. Whether autonomy is a diminution of autonomy depends upon whether a person is **intentionally** directed to the point of determinations and acts. In other words, whether autonomy is compromised depends on whether a person is under the control of others. The others, here, are those who are intentionally going to govern his/her decisions and acts. From a viewpoint of medical advice of healthcare providers, if a patient is operated upon or deceived by them and makes a decision unconsciously which they intend, a patient's autonomy will have been diminished with regard to his/her medical decision-making.

4. No compromising a person's autonomy

If the intentional control by others is not imposed, the autonomy of the person is not compromised. The person still directs their own decisions and acts and can be fully autonomous with regard to them. In that case, the degree of achievement of his/her purpose need not concern us, and it can be said that he/she is autonomous.

5. "Instrumental value of autonomy" is the ethical foundation of informed consent

When a person is not controlled intentionally, and when his/her aims are not achieved by lack of knowledge, **the autonomy itself** of the person was not compromised and only **the instrumental value of autonomy** of the person is compromised. As for the autonomy of the patient with regard to the treatment decision, it will not be compromised even if healthcare providers failed to secure the informed consent of the patient as a result of having non-intentionally or carelessly neglected to disclose relevant information.

6. Achievement of value and the purpose -Securing of "instrumental value of autonomy"

In this way, respect for autonomy will not be the only theoretical basis for securing informed consent. It can, nonetheless, be fair grounds to secure informed consent if we think achieving respect for autonomy secures the instrumental value of the autonomy. As we argued in the beginning, healthcare providers should support patients in such a way that patients can achieve their purposes in medical care. All we need for this is dialog and participation of patients and healthcare providers.

7. Shared decision making for patient's best interests

In order to respect a patient's autonomous choice and for a patient to make the final choices, healthcare providers need to know a patient's purposes and values and to offer information required for medical treatment choices. For that purpose, healthcare providers should give patients information required for choice; patients should reflect on their purposes and values, and both should have a mutual dialog for the purpose of knowledge acquisition regarding the medical treatment. That is to say, in order to secure suitable informed consent, the dialog, it is indispensable to have a sufficient communication about medical treatment between healthcare providers and patients for achieving agreement and selection, and participation of both in decision-making. The patients can, in this way, make autonomous decision that will satisfy their own purposes and values.

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The concept of objectification in life philosophy and natural science

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Introduction

In this presentation, I will focus on two German philosophers, who dealt with problems of human life and attempted to establish a rigorous method equivalent to that of natural science. My aim here is to elucidate the differences and similarities in the meaning of the term '*objectification*' as used in life philosophy and natural science.

1. Life philosophy and Heidegger's early thought

During the time of his early studies in Freiburg (1919-23), Martin Heidegger analyzed human life in a philosophical way. In that period, *Life philosophy* still held a strong influence in Continental philosophy. At that time, with the prosperity of natural science, application of positivistic approaches to humanity had become popular. Against such currents, life philosophy generally attempted to illuminate human life in a different way to experimental psychology or biology. This movement of thought, which reflected the situation of the period, inspired the young Heidegger's ideas.

Heidegger appreciated the effort and consciousness of problems in life philosophy. He often mentioned it in his early lectures. Wilhelm Dilthey had an especially profound influence on Heidegger in forming a perspective from which to approach human life. Heidegger shared the motivation of Dilthey to search for life's strict expression in terms other than those of natural science.

2. Differences in attitude toward scientific objectification between Heidegger and Dilthey

The two philosophers take contrasting attitudes, however, toward the objectification of natural science. Dilthey, especially in the autumn of his life, wrestled with life's objective expression. He tried to express the inner structure of life objectively in the relation to the historical world. His final aim was to make a science of human life which bears comparison with natural science.

For the purpose of research into human life he avoided the metaphysical route and focused somewhat on evaluation of the empirical method. This was because it was important for him not to transcend but to base his thinking on each actual experience. So he appreciated the own value of natural science and adopted the criterion of scientific objectification, which is formed by generalization.

On the contrary, Heidegger rejected scientific objectification radically. He thought that human life is never expressed in its own structure in such a way. So he tried to generate the original method with which to approach human life not relying on the scientific objectification.

3. Dilthey's project of human studies

Dilthey is known as a founder of human studies (*Geisteswissenschaften*) as a domain independent of natural science. According to him, human beings are both natural and historical beings. And the object of human studies is historical human life, which is never grasped by the scientific *explanation*. The historical human life is only approached by *understanding*. Understanding never aims to build universal principles like natural science. It aims rather to describe the individual, which always takes its own position in each historical context.

At the same time Dilthey attempted to make a bridge between a human individual life and objective spirit. He thought that every individual life is comprehended in a common historical world. Furthermore, each individual receives each meaning of presence from this world. He also thought that a finite individual life can express and understand itself objectively or universally in a similar way to natural science by way of human studies. So according to him natural science and human studies are concerned with two different realms of life, but both are equally strict sciences.

In this way, although he clearly distinguished the methods of natural science and human studies, eventually he relied on scientific objectivity, namely, universal and valid objectivity, to secure the academic rigor of human studies. Here we can find the great difference between Dilthey and Heidegger.

4. The genesis of the scientific objectification

More radically than Dilthey, Heidegger thought that human lives can never be expressed *objectively* in the meaning of natural science. Because he thought that by the process of such objectification life is separated from itself (*Entlebung*). He thought the genesis of the scientific objectification as figured below.

Life's daily region

- *instability
- *under the control of the accidental event or each individual tendency
- *continual deficiency and complacency
- *uncertainty

*Tendency to stabilize and get the possibility of prediction



*Tendency to generalize and objectify
tendency to build universal principals



*Readiness for the occurrence of each science

According to Heidegger, life's realm is always exposed to instability and uncertainty. Its instance is for example unforeseen accident and especially each person's death. He thought that such unstable characteristic of life made ready for scientific attitudes of objectification. However, he thought that life itself would never be understood in this way like other particular sciences, which begin with the tendency to stabilize. Life itself is essentially uncertain and instable. Heidegger called the tendency to stabilize *the fall* (*Ruinanz*). And in order to trace back to the origin which doesn't suffer the tendency to objectify, philosophy of life must proceed against this tendency.

5. Objectification to which philosophy should aim

But on the other hand, philosophy as one science (*als Wissenschaft*) can never be completely separated from *objectification*. It needs a strict formation of concepts. So Heidegger attempted to form concepts without injuring the activity or historicity of life. Heidegger didn't reject the objectification itself. He tried, however, to regenerate the meaning of objectification which never depends on scientific meaning.

For Dilthey, the main problem of human studies was describing the individual. Therefore each individual needs to overcome its finiteness of individuality in order to express and understand itself not arbitrary but objectively. For that purpose the existence of the human history, community, or culture that incorporates the individual were very important. Understanding the individual experience can be accomplished sufficiently only through the historical world which is common and open for all the human spirit.

On the contrary, Heidegger thought that the finiteness of life can never be overcome and rather philosophy should go along with this finiteness. He became conscious of a deficiency in describing the individual life as Dilthey had done. Expression by describing has always a tendency to separate from the finite life and to tend toward the infinite.

According to Heidegger, philosophical objectivity means "*not being apart from the acting life*". Philosophical approaches to human life should always be coordinated with each irreplaceable acting life. The meaning or criterion of objectivity in philosophy is never fixed and rather it should be constantly renewed. For that purpose the objectification he aimed to has always a characteristic of incompleteness.

In his early lectures, Heidegger mainly tried to disclose the wrong application of objectification to human life. So in his descriptions negative expressions appeared frequently. He was in that time on the way to discover his own philosophical attitude different from the methods of natural science and also different from that of Dilthey. And eventually in "Being and time" (1927), he presented first his own systematical construction of concepts. Then his target of consideration diverted from life to *being*. His research into strict methodologies led him to investigate the meaning of *being*. This transition of the theme occurred through his attitude toward philosophical objectification, which continually requires renewal to be along with acting life.

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On the Myth-Making Function in Bergson

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Introduction

Purpose

To show clearly the originality of 'the myth-making function' (la fonction fabulatrice) in Henri Bergson through a comparison with 'the creative imagination' (l'imagination créatrice) in Théodule Ribot

Background

- According to Bergson, the myth-making function is the virtual instinct in human being that restrains intelligence from excessive reflection.
- Several studies on the myth-making function have focused on the relation between instinct and intelligence.
- Bergson also insists that the myth-making function is a sort of imagination, but yet there is a difference between the function and psychological imagination.

Methods & Materials

Methods

- To concentrate on following issues:
 - 1. imagination and intelligence in Bergson and Ribot**
 - 2. imagination and emotion in Bergson and Ribot**
- To analyze both descriptions about the relation between imagination and intelligence, and also imagination and emotion

Materials

- Intelligence is guided in fact by present perceptions or by that more or less vivid residue of perception called recollection (MR, 126).
- it [a virtuality of instinct, the residue of instinct] cannot exercise direct action, but, since intelligence works on representations, it will call up "imaginary" ones, which will hold their own against the representation of reality and will succeed, through the agency of intelligence itself, in counteracting the work of intelligence (MR, 124).
- the emotion felt by a man in the presence of nature certainly counts for something in the origin of religion. But, we repeat, religion is less a fear than a reaction against fear, and it is not, in its beginnings, a belief in deities (MR, 160).
- this demand [a fundamental demand of life] has called into being the myth-making faculty; the myth-making function is thus to be deduced from the conditions of existence of the human species (MR, 207).
- the creative imagination has the faculty of *thinking analogically* [...] as an essential and fundamental element in the order of intelligence (IC, 22).
- The imagination in the order of intelligence is equivalent to the will in the order of action (IC, 6).
- the creative imagination implicates affective elements [, and] all inventions presuppose a demand, a desire, a tendency, an impulse (IC, 27).
- All the creations are useful practically and a demand for life has produced them (IC, 38).

Conclusion

Conclusion

The myth-making function is

- (1) the creative function that makes representations; therefore the myth-making function can take the place of intelligence. However, the same thing may be said of the creative imagination.
- (2) the function that is made by *the emotion*, in other words, **the demand of life makes the myth-making function**. On the other hand, the creative imagination is not made by any emotions, but only may be affected by emotions.

Future Discussions

Why can the human will to live outwits *the nature* that made insects, animals, human beings and so on, by means of the myth-making function or intelligence?

Results & Discussions

Results & Discussions

- Both the myth-making function and the creative imagination make representations that look just like present perceptions to disturb excessive intellection, being affected with emotions (for example, fear, pity, joy and so on).
- Bergson divides the emotions into two groups:
 - (1) infra-intellectual emotion**
 - (2) supra-intellectual emotion**
- In Bergson's view, in fact, there are some emotions that are excited by perception and representation. Generally speaking, psychologists discuss these emotions **(1)**.
- However, there is also the emotion as the demand of life that forms the basis of all the creation and stimulates the myth-making function to make representations, images, mythology, religious ceremonies, literatures, etc **(2)**.
- The myth-making function and the creative imagination are more or less similar, but **Bergson clarifies not only the source of imaginative creations but also the source of imagination (the myth-making function)**. In this sense, the myth-making function is distinct from the creative imagination.

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'Prolegomena' Revisited

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Introduction

Husserl presented the structure of sciences as systems of right knowledge in the 'Prolegomena' in his *Logical Investigations*. That is what he calls "pure logic". Sciences have the character of a system of knowledge precisely because they contain pure logic in themselves. To defend pure logic, in 'Prolegomena' Husserl denies psychological interpretations of logic, relativistic arguments to truth, or the skepticism to right knowledge. The essential argument among these arguments is the criticism of Protagorean skepticism; Husserl thinks that every relativism or skepticism has the Protagorean view in common.

Protagorean skepticism and Husserl's criticism of it

Protagorean skepticism claims that the individual man is the measure of all things. "For each man that is true which seems to *him* true, one thing to one man and the opposite to another, if that is how he sees it"(LI). Thus "all truth (and knowledge) is relative - relative to the contingently judging subject". Hence Protagorean skepticism denies general truth or truth in itself. Husserl refutes Protagorean skepticism "in an objectively valid manner": "laws such as the law of contradiction have their roots in the mere meaning of truth", and "from these it follows that talk of subjective truth, that is one thing for one man and the opposite for another, must count as the purest nonsense". This is also because "in setting up his theory he (the skeptic) is making a claim to be convincing to others, a claim presupposing that very objectivity of truth which his thesis denies".

The first problem of the criticism

Husserl's criticism is relevant to a kind of Protagorean skepticism. In this kind of skepticism the skeptic himself knows what so called truth is or what right knowledge is. Otherwise he could not refute the other theories in general and set his own theory as true.

Can Husserl's criticism eliminate all kinds of skepticisms? Could not there remain any kind of skepticism? It seems that a skepticism which can avoid this criticism remains. That is Pyrrhonian skepticism. What I call here Pyrrhonian skepticism is a theory which claims that "if people who hold beliefs posit as real the things they hold beliefs about, while Skeptics(Pyrrhonian skeptics) utter their own phrases in such a way that they are implicitly cancelled by themselves, then they cannot be said to hold beliefs in uttering them"(OP), and hence that "in uttering these phrases they(Pyrrhonian skeptics) say what is apparent to themselves and report their own feelings without holding opinions". This skepticism implicitly shows that it is insufficient to draw a mere figure positing ego and objects or contents. The very functions of various assessments (doubtful, sure and so on) of objects or contents by ego must be considered. Then truth or knowledge cannot be defended unless Pyrrhonian skepticism, which does not decide whether or not truth in itself exists is refuted. Thus the task of refuting this type of skepticism still remains. This is the first problem which Husserl's criticism involves.

The second problem of the criticism

It follows that Husserl so to speak naively posits being-in-itself, such as truth in itself because he doesn't focus on acts which assess ontological entities of objects in the "Prolegomena"; this is the second problem of Husserl's criticism. When skepticism is given right consideration, it reveals various acts of comprehending objects by ego. Indeed, skepticism denies truth and knowledge. Therefore it is essential to affirm general truth in itself to defend sciences and knowledge. But is it not until skeptical problems are addressed rightly that we are able to identify and approach general, epistemological problems about subject and its acts?

In this sense, it is interesting that Husserl consistently and naively posits being-in-itself, such as truth in itself in *Logical Investigations*. It is true that the way that how being-in-itself exist is clearly distinguished from how real things exist, but when someone asks "so, where and how is that being-in-itself", there is no choice but to answer "somewhere and somehow" as far as *Logical Investigations* goes. Another detailed consideration about ontological entities of being-in-itself is necessary.

This does not necessarily mean that the arguments about epistemology in *Logical Investigation* are naive. As we know, Husserl addressed epistemological problems including how we acquire the general and what intentionality is in "Investigation II", "Investigation V", and "Investigation VI". Though these investigations are thorough; being-in-itself remains ambiguous. This is partly because Husserl did not focus on elements of ego in considerations of the "Prolegomena". This is a pity, though it is all the better for frequent criticisms of Husserl's theory about being-in-itself. If Husserl had already contemplated various assessments of objects by ego in *Logical Investigations*, the arguments about ontological entities of being-in-itself would be more interesting.

Conclusion

In summary, Husserl's criticism of skepticism can be summarized as follows: Husserl posits truth in itself because he claims all theories about logic presuppose truth or knowledge and he also defends them, but there is room for doubt that truth in itself exists in a rigorous sense because his argument in the 'Prolegomena' does not consider Pyrrhonian skepticism. On top of that being-in-itself is left ambiguous because Husserl overlooked the problem of assessment acts by ego. It might possibly seem that the way in which he addresses the problems from the 'Prolegomena' is non sense, but we must not evaluate his arguments such a way. This is because they are meaningful when they are considered in the course of the development of Husserl's thought. In fact we should contemplate even obvious faults in his arguments by contrast with Husserl's later thought. This task is left to contemporary researchers.

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The problem of the relationship between individuality and universality in Hegel's philosophy

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Introduction

One of the most important motives in Hegel's philosophy is that of the subjectivity. This motif is derived from the argument over the relationship between individuality and universality in the medieval times. — It is further formed by Hegel's attempt to respond to the problem of the character and status of cognition, which has been one of the matters of concern for Modern philosophy since Descartes.

1.

It is characteristic of Hegel, that the problem of individuality and universality receives treatment as “-logical-” in Hegel's sense of the word, that is to say it is treated as a -metaphysical problem-. The common sense and common place understanding of “individuality” is that it is a word that signifies this or that particular one who recognizes or thinks. Hegel's understanding however is different. According to his theory, individuality and universality are not different from each other originally.

His thought is based on the notion, that the universality is a movement which develops itself. Universality has, unlike Aristotelian substance (οὐσία), not stationary attributes. It forms itself however into one conclusive process by determining itself through its ownself, that is, by particularizing itself. — In this way the character of the universality as subjectivity consists in developing a process. Hegel recognizes the individuality as the subjectivity in the one process, for this is the conclusion which is drawn from the nature of the subjectivity.

2.

It is significant that the aim and end of the process is nothing other than this subjectivity itself. The subjectivity arrives back upon itself through the movement of determination and particularization. This movement is called “Reflexion”. — The notion of arriving back implies that the subjectivity, which goes through the process, is by nature that, which it has been. — But “by nature” means by no means that it has been so before attributes have been predicated in the case of substance. — “That what it has been by nature” however reveals itself with the principle “to know” —; and it is nothing else but “to know”.

Therefore this self-developing-process of subjectivity functions as a bringing out of the internal structure of self-consciousness, which is the most important motif in Modern philosophy.

We are apt to delude ourselves into believing that the self-consciousness is self-evident to anyone. it is not easy, however, to state what it is in its essence. To solve this difficulty, Hegel provides an answer with his theory of the self-development of subjectivity as arriving back upon itself. — This is also expressed as “the Concept is the personality”. The term “concept” is the key-word not only for his “logic”, but for his whole system of philosophy, permutes his notion of subjectivity which itself contains universality, individuality and particularity as its element.

3.

The notion “personality” as a philosophical term has been formed by degrees through efforts to understand the meaning of the Holy Trinity (Father, Son and the Holy Spirit) in Christian theology. — St. Augustine of Hippo (354~430) for example, explains the notion of “—one essence and three persons—” by means of triplicity of the mind (memory, intelligence and volition). —

Nicholas of Cusa (1401~1464) is the one who opens the way for Hegel's theory of subjectivity. — In the one of his chief works “ —On no other— ”, Nicholas develops his theory of self-determination of the universality (—i.e. God) with the notion of the “ —no other— ”. — Everything that exists is, so far as it exists, determined. The cognition, which recognizes it is also determined. For existence and cognition the “ —no other— ” is their principle and therefore not differs from themselves. — From this viewpoint Nicholas understands the individuality and universality as a unity.

In the foundation of Hegel's subjectivity theory Nicholas's thought plays an important role in establishing the notion that the universality is a movement which develops itself. As the background of Nicholas's “ no other ” plays the theological motif creation. It is however not to disregard that Nicholas would show the creator and the creature not as two quite different things but they are in unity. Creation is unlike this or that act which is done commonly, by no means any particular performance at this or that point in time. It means to bring the creature into that what it has been. Nicholas learns this thought by the Latin commentaries of the Bible written by Meister Eckhart.

We should not indeed forget that Nicholas and Eckhart's thought have a theological origin. We must however on the other hand take their incomparable and so profound thought in order to understand Hegel's subjectivity and personality of the concept more deeply and radically. And we can gain an wider perspective for the subjectivity, which is one of the most important motives of Modern philosophy and metaphysics in general. The problem of subjectivity is a problem of metaphysics which requires further examination from more far-reaching perspective.

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Hume on Logic and Demonstration

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Introduction

As is commonly known, the 18th century Scottish philosopher David Hume presented seven philosophical relations in his *Treatise*. Later these were divided into two subjects as all objects of human reasoning or inquiry in his *Enquiry*; 'relations of ideas' and 'matters of fact'. The former are objects of our knowledge and certainty. These were resemblance, contrariety, degrees in quality, and proportions in quantity and number. According to him, these relations are intuitively certain, except equality or any exact proportion in the relations of proportions in quantity or number. Equality or any exact proportion is demonstrably certain. Besides the question of what intuition might be as an operation of our minds, demonstrations relating to intuition in Hume's arguments are significant. Because Hume was primarily concerned with revealing faculties or operations of our mind: human understanding. On the other hand, he definitely stated in his Abstract that 'the author[Hume himself, in fact] has finished what regards logic[in Book I of *Treatise*]', impersonating a book reviewer of it. However his arguments on human understanding doesn't surely seem to be engaged in logic at least in our time – modern formal logic. Thus Hume's logic must be considered as different mode of logic from our logic or Aristotelian school logic.

1. The Early Stage of Hume's Logic

(1) Traditional Logic in School

Aristotle's *Organon* was a core text of European universities during the 14th and 15th centuries. None the less, his *Topics*, and the like, works that take up the ways or laws by which we practically reason in various situations was rarely included in university curricula of that time. Though unlike Russell said so, Aristotle himself by no means thought little of inductive reasoning, medieval logicians, philosophers, and theologians placed high value exclusively on deductive arguments that were constructed from categorical propositions.



(2) The Heuristic Methodology

The flourishing of natural philosophy in early-modern times originates from Copernicus, and is chiefly due to the adoption of experimental methods by practitioners of natural philosophy outside the universities. This method in their inquiry depended not on deductive arguments but on inductive reasoning, or reasoning from cause to effect, as a principle of thinking. For new natural philosophers, unlike neo-Aristotelian natural philosophers, who had set the notions of *generalized experiences*, had to appeal to the particular experiences of natural phenomena for seeking out *new* phenomena of nature.



(3) The Bridge between Natural Philosophy and Hume

For Francis Bacon, whom Hume considered as 'the father of experimental physics', induction is logic itself and also scientific method. He said, "we reject proof by syllogism". The reason Hume appreciated him, however, was related I think to Bacon's approach to the way our reasoning or thought should be considered together with faculties or operations of our mind. Thus Bacon is an originator, who began "to put the science of man on a new footing". Later Hume attempted to establish 'the science of man' of his own, with 'the experimental method of reasoning' from Newton's 'Rules for the study of natural philosophy' in *Principia* - *inductive reasoning, or reasoning from cause to effect*.



*Comment:
According to Gaukroger, the period between the beginning of 16th century and the publication of Boole's *Mathematical Analysis of Logic* (1847) is regarded as 'the interregnum in the history of logic'(Gaukroger intro 1).

2. The Naturalistic View of Logic

(1) Hume's View of Logic

So, what does logic mean for Hume? I think that to answer that question, we should pay attention to Hume's criticism on existing logic that was taught at schools in those days. Though Hume restrained related the criticism in a footnote to the *Treatise*, it is significant to an understanding of his view of logic.

"This error consists in the vulgar division of the acts of the understanding, into *conception, judgment, and reasoning*, and in the definitions we give of them. [...] But these distinctions and definitions are faulty in very considerable articles. [...] What we may in general affirm concerning these three acts of the understanding is, that taking them in a proper light, they all resolve themselves into the first, and are nothing but particular ways of conceiving our objects. [...] the act of the mind exceeds not a simple conception". (THN 1.3.7.5.n2)

Leibnitz's Indication of the Common Systems of Logic toward the Copiousness in Explaining the Operations of the Understanding (ATHN 4)

(2) Arnauld & Nicoles' View of Logic

There is a text book of logic, A. Arnauld & P. Nicole's *La Logique ou l'Art de Penser*, that Hume surely read, and which he would apparently (though partly) see as a target of his criticism on school logic. According to this book, "logic is the art of conducting reason well in knowing things, as much to instruct ourselves about them as to instruct others", and the art consists in *reflections*, which are made four operations of our mind: *concevoir, juger, raisonner, and ordonner* (Arnauld & Nicole 23).

However, there are plausible reasons for Arnauld & Nicole to have enumerated the operations of our mind. They aimed at educational effectiveness in presenting their material to those, who had not studied any rules or disciplines of logic. Moreover they thought that operations of our mind in *reasoning are done naturally* (ibid.).

3. A Different View of the Same Object

(1) Two Definitions of Cause

Hume presented famous two definitions of cause in his *Treatise* and *Enquiry*. It is commonly said that, one is from the point of philosophical relation; the other is from the view of natural relation (THN 1.3.14.30; EHU 7.2.4).

D1: "An object precedent and contiguous to another, and where all the objects resembling the former are plac'd in like relations of precedence and contiguity to those objects, that resemble the latter".

D2: "A CAUSE is an object precedent and contiguous to another, and so united with it, that the idea of the one determines the mind to form the idea of the other, and the impression of the one to form a more lively idea of the other".

↓
In these definitions, it is essential to count them as drawn from the objects whether foreign to cause or domestic to cause. The former is independent of faculties or operations of our mind, while the latter depends on faculties or operations of our mind, that is, on the principle of '*associations of ideas*', in spite of our perceiving the same objects.

(2) General Rules by which to Judge of Cause & Effect

On the other side, Hume presented "general rules, by which we may know when they [cause & effect] really are so". From first to third rules are articulated about 'contiguity of time and space of cause and effect', 'precedence of cause', and 'constant conjunction of cause and effect', fourth rule confirms the principle of uniformity of nature, and so on.

↓
After fixing these rules provisionally, he claimed that "here is all the Logic I think proper to employ in my reasoning". Moreover he commented on them as that "perhaps even this was not very necessary, but might have been supply'd by the natural principles of our understanding". (THN 1.3.15.3).

Hume's Criticism of a Long System of Rules and Precept by Scholastic Headpieces & Logicians (THN 1.3.15.3)

Conclusion

- (1) For Hume, logic is a set of provisional and elaborate rules to deal with unexperienced particular object, with the method of inductive reasoning or reasoning from cause to effect, to be aimed for a higher level of assurance, which could be strengthened by the evidences of repeated experiences and general or historical confirmations. Demonstration is simply our discovering the relations of ideas.
- (2) If this is ever appropriate understanding of Hume's argument of logic, he could be practically on a par with the claim that all the reasoning we employ would result in probable reasoning, without appealing to a certain exemplar model, including proportions in number in philosophical relation.
- (3) Historically speaking, he argued against the traditional scholastic doctrines of logic that had made unreflectingly or dogmatically represent *terminus* something ambiguous as materials of logic.

4. Hume on Demonstration as a Philosophical Relation

Let's see what Hume calls demonstration. We should notice that these four relations, including latter three relations, are presented as philosophical relation, which is applied for "particular circumstance, in which, even upon the arbitrary union of two ideas in the fancy, we may think proper to compare them" (THN 1.1.5.1).

As I said earlier, only in cases of equality or any exact proportion in the relations of proportions in quantity or number, we can determine their relations without any error. For when we do so, we are possessed of a precise standard (i.e. a unit) to judge them. But unfortunately Hume didn't exhibit an clear example of demonstrative reasoning. D. C. Stove (or Broughton, Mackie, Beauchamp & Rosenberg), for example, claimed that what Hume meant by demonstrative arguments was strictly deductively valid arguments with necessarily true or a priori true premises (Stove 35-36). It would be one of reasons that skepticism about induction and logical necessity about deduction are like two sides of the same coin in interpretation on Hume. On the other hand, D. Owen (or Beebe) interprets Hume's original form in philosophical relation of demonstration as follows (Owen 94-96).

An Example of Demonstration of the Equality of " $5^3 + 12 = 137$ "

- ① We begin with intuiting that the idea of 5^2 and the idea of 25 stand in the relation of equality each other.
- ② At the same time, we intuit the proposition of $5 \times 5 = 25$.
- ③ Then we intuit the proposition of $5^2 \times 5 = 125$.
- ④ And then we intuit the proposition of $5^3 + 12 = 12 + 125$.
- ⑤ Finally we intuit the proposition of $5^3 + 12 = 137$.

**Comment:
At least, the term 'deduction' was not used for limited expression of deductively valid arguments with necessarily true or a priori true premises in that days (Volton 92-93; Owen 117, et al.). Hume actually used the word 'deduction' as exclusively 'reasoning' or 'argument'.

I agree basically with Owen. To be sure, he never rejected demonstrations *in itself* as the laws of thinking. But that derivation was not the principal interest of Hume. Because he gave a thought to this problem as follows.

"Truth is of two kinds, consisting either in the discovery of the proportions of ideas, consider'd as such, or in the conformity of our ideas of objects to their real existence". (THN 2.3.10.2)

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The rationality of the Precautionary principle: making the precautionary principle more applicable

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Introduction

Scientific uncertainty

We must make decisions in circumstances in which insufficient knowledge of possible high risk events compromises the reliability of scientific analysis.

The rationality of the Precautionary principle

The Precautionary Principle (PP) requires us to adopt approaches in order to protect human health and the environment despite a lack of sufficient scientific certainty.

Some doubt the applicability of PP, claiming that it lacks a rational basis and may thus lead to irrational conclusions.

This might suggest that we cannot make public policy decisions by means of the precautionary principle alone.

I would like to examine the notion of “rationality” and to emphasize the implications of “deliberative democracy” with regard to the better application of PP.

Risk analysis and management

Risk-Cost-Benefit Analysis

There exists an important question: introducing the new technology, can we make profits commensurate with costs and risks?

We need “risk-cost-benefit analysis(RCBA)”: the method of calculating cost of the introduction of new technology and predicting risks and benefits.

Problems of risk management

RCBA (risk-cost-benefit analysis): Such procedures embody uncertainties; models used in RCBA are not always perfectly sophisticated models of real events.

In principle, we never entirely know what happened, happens and will happen in the real world, which makes us more modest about our actions that might pose harm to human health or the environment.

Scientific Uncertainty

The limits of our knowledge and science

As science is the process for exploring veiled realms, in cutting-edge research scientific knowledge is constantly renewed by new scientific discoveries.

There are, however, some cases where we cannot elucidate the causal mechanism of complex phenomena at the present time due to the complexity of phenomena (ex. global warming)

Thus, scientific uncertainty is due to the plasticity or limits of our knowledge.

Attitudes towards scientific uncertainty

In regulating or prohibiting a certain product or procedure, it seems to be necessary to prove the causal relationship between products and the possible harmful effect.

As we know, since scientific knowledge might be renewed, we must not have too much confidence in science.

Therefore, we, especially policy-makers, must keep in mind that science has limitations.

Precautionary Principle

The Emergence of PP

PP has been adopted in the legislation for protection of human health or environment.

The most prominent and frequently cited version of PP is probably the 1992 *Rio declaration*:

“In order to protect the environment, the precautionary approach shall be widely applied by States according to their capabilities. Where there are threats of serious or irreversible damage, lack of full scientific certainty shall not be used as a reason for postponing cost-effective measures to prevent environmental degradation.”

One of the most remarkable characteristics of PP is that even if there exist no full scientific evidences it is possible and legitimate to regulate or prohibit products or procedures which might pose serious and irreversible damage to human health or the environment.

The defects of PP

Opponents characterize the definition of PP as too vague to serve as a regulatory standard.

Because of its vagueness PP might lack rational foundation as the principle for action.

Is PP rational?

Perhaps its vagueness let us use PP arbitrarily; the government use PP in order to defend the domestic market from imported goods.

It might be impossible that we lead to rational consequences by means of PP alone.

Therefore, we have to consider the implication of PP as a norm in the context of society.

The future generations

Effects on the future generations

In the era of modern science and technology there exist threats to health or the environment of the future generation(ex. nuclear waste, greenhouse gas effects etc.)

Must we pay homage to the rights of the future generation which doesn't exist at present?

Moral obligations to the future generations

Intergenerational ethics: requires us to keep in mind the benefits of the future generations.

We have to preserve the living conditions of the future so as not to damage their welfare.

Deliberative Democracy

Social justice and democratic legitimacy

There is an ideal that all participants in public decision making should be well informed and in dialogue with each other voluntarily.

This ideal is based on the belief that reasonable consensus should be derived from dialogues.

Deliberative democracy and PP

Precautionary measures are fundamental elements of the ideal democratic decision making of risks imposed on the future generations by the present generation.

We must pay homage to the future's ability to make decisions in order for them to be informed well and in dialogue voluntarily.

Conclusions

- We must perform risk management on the grounds that there might be scientific uncertainty.
- Without support of social consensus, the PP would not function well.
- We must maintain conditions of living by which the future generations can live, as well as our own conditions.
- In deliberative democracy we can take precautionary measures to make sustainable developments possible.

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Viscosity solutions on a Riemannian manifold

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Preliminaries

Level Set:

For an almost everywhere differentiable function $u : \mathbb{R}^N \times [0, T] \rightarrow \mathbb{R}$ ($\nabla u \neq 0$ a.e.), define level set at time t as,

$$\Gamma_t = \{x \in \mathbb{R}^N : u(t, x) = 0\},$$

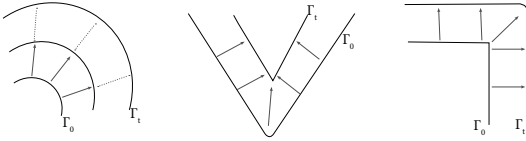


Figure 1. Evolving surfaces with their unit normal vectors.

Surface Evolution Equation, In general

$$V = G(t, x, \mathbf{n}, \nabla \mathbf{n}) \dots \dots \rightarrow$$

for $t > 0$, and $x \in \Gamma_t$

V : normal velocity ($:= \frac{u_t}{|\nabla u|}$)

\mathbf{n} : unit normal vector ($:= \frac{\nabla u}{|\nabla u|}$)

using level set representation

$$u_t + F(t, x, u, \nabla u, \nabla^2 u) = 0$$

where,

$F : (0, T) \times \Omega \times \mathbb{R} \times \mathbb{R}^n \times \mathcal{S}^n \rightarrow \mathbb{R}$
 \mathcal{S}^n , space of symmetric matrices.

Example:

1. $V = H$ (Mean curvature flow) $\Rightarrow u_t + |\nabla u| \operatorname{div} \left(\frac{\nabla u}{|\nabla u|} \right) = 0$
2. $V = K$ (Gaussian curvature flow)

Semijets:

Let $f : (0, T) \times M \rightarrow \mathbb{R}$ be a lower semicontinuous function (upper semicontinuous function). Define a parabolic second order **subjet** (**superjet**) of f at a point $(t_0, x_0) \in (0, T) \times M$ by

$$\mathcal{P}^{2,-} f(t_0, x_0) (\mathcal{P}^{2,+} f) := \left\{ (D_t \varphi(t_0, x_0), D_x \varphi(t_0, x_0), D_x^2 \varphi(t_0, x_0)) : \varphi \in C^{1,2}((0, T) \times M) \text{ and } f - \varphi \text{ attains a local minimum (max.) at } (t_0, x_0) \right\}$$

where, by $\varphi \in C^{1,2}((0, T) \times M)$, we mean that φ is once continuously differentiable in $t \in (0, T)$, and twice continuously differentiable in $x \in M$.

Here, $\mathcal{P}^{2,-} f(t, x)$, and $\mathcal{P}^{2,+} f(t, x)$ are subsets of $\mathbb{R} \times T_x M^* \times \mathcal{L}_s^2(T_x M)$, where $T_x M^*$ is the cotangent space at $x \in M$, and $\mathcal{L}_s^2(T_x M)$ denote the space of bilinear forms on $T_x M$, and φ is defined in a neighbourhood of (t_0, x_0) .

Cauchy-Dirichlet problem:

We consider the following Cauchy-Dirichlet problem of the form:

$$[CDP] := \begin{cases} u_t + F(t, x, u, Du, D^2 u) = 0 \text{ in } (0, T) \times \Omega \\ u(t, x) = h(t, x), (t, x) \in [0, T] \times \partial\Omega, \\ u(0, x) = \psi(x), x \in \bar{\Omega}, \end{cases}$$

where u is the function of $(t, x) : [0, T] \times M \rightarrow \mathbb{R}$ and M is a finite-dimensional complete Riemannian manifold, $Du, D^2 u$ mean $D_x u(t, x), D_x^2 u(t, x)$, Ω is open and bounded in $M, T > 0, h \in C^{1,2}([0, T] \times \bar{\Omega})$ and $\psi \in C^2(\bar{\Omega})$.

Definition(Viscosity solution)

Let M be a Riemannian manifold, and $F \in C(\mathcal{X}, \mathbb{R})$. We say that an upper semicontinuous function (lower semicontinuous function) $u : [0, T] \times \bar{\Omega} \subset M \rightarrow \mathbb{R}$ is a **viscosity subsolution** (**viscosity supersolution**) of [CDP] on $[0, T] \times \bar{\Omega}$, if

$$\begin{cases} u(t, x) \leq (\geq) h(t, x), (t, x) \in [0, T] \times \partial\Omega, \\ u(0, x) \leq (\geq) \psi(x), x \in \bar{\Omega}, \end{cases}$$

hold for all $(t, x) \in (0, T) \times \Omega$, and for $(p, \zeta, A) \in \mathcal{P}^{2,+} u(t, x)$ ($\in \mathcal{P}^{2,-} u(t, x)$),

$$p + F(t, x, u(t, x), \zeta, A) \leq 0 (\geq 0),$$

hold.

If u is both a viscosity subsolution and a viscosity supersolution of [CDP], we say that u is a **viscosity solution** of [CDP].

Here,

$$\mathcal{X} := \left\{ (t, x, r, \zeta, A) : t \in [0, T], x \in M, r \in \mathbb{R}, \zeta \in T_x M^*, A \in \mathcal{L}_s^2(T_x M) \right\}.$$

Key Tools

Properness of F :

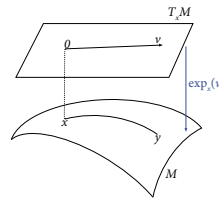
We will say that a function $F : \mathcal{X} \rightarrow \mathbb{R}$ is **proper** if satisfies;

(i) F is degenerate elliptic provided that

$$P \leq Q \Rightarrow F(x, r, \zeta, Q) \leq F(x, r, \zeta, P),$$

(ii) F is nondecreasing in the variable r .

Exponential map:



Jacobi field:

Let $\gamma : [0, l] \rightarrow M$ be geodesic, and $\alpha(t, s) : [0, l] \times [-\varepsilon, \varepsilon] \rightarrow M$ variation through geodesics $\alpha(t, s) = \gamma(t)$. Any vector field satisfying the Jacobi equation

$$\left(\frac{\nabla}{dt} \right)^2 V = R(\gamma'(t), V(t)) \gamma'(t)$$

along a geodesic is called a **Jacobi field**.

Proposition:

For the function $\varphi(x, y) = d(x, y)^2$ on $M \times M$,

(1) If M has nonnegative sectional curvature then

$$\langle D^2 \varphi(x, y)(v, L_{xy} v), (v, L_{xy} v) \rangle \leq 0.$$

(2) If M has nonpositive sectional curvature then

$$\langle D^2 \varphi(x, y)(v, L_{xy} v), (v, L_{xy} v) \rangle \geq 0.$$

Here, we denote by L_{xy} the parallel transport along the unique minimizing geodesic connecting x to y

Theorem(Comparison Principle)

Let Ω be a bounded open subset of a complete finite-dimensional Riemannian manifold M with nonnegative sectional curvature, and $F : \mathcal{X} \rightarrow \mathbb{R}$ be proper for each fixed $t \in (0, T)$ and intrinsically uniformly continuous with respect to x , uniformly in t . Let $u \in USC([0, T] \times \bar{\Omega})$ be a subsolution and $v \in LSC([0, T] \times \bar{\Omega})$ a supersolution of [CDP]. Then $u \leq v$ on $[0, T] \times \Omega$.

M has nonnegative sectional curvature

\Downarrow implies

$P \leq L_{yx}(Q)$ and degenerate ellipticity

\Downarrow implies

$$F(t, y, r, L_{xy} \zeta, Q) - F(t, x, r, \zeta, L_{yx} Q) \leq w(\alpha d(x, y)^2 + d(x, y)).$$

\Downarrow means

F is intrinsically uniformly continuous.

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Smoothness of densities of generalized locally non-degenerate Wiener functionals

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1 Introduction

Several criteria for existence of smooth densities of Wiener functionals are known in the framework of Malliavin calculus. In this presentation, we introduce the notion of generalized locally non-degenerate Wiener functionals, and prove that they possess smooth densities. The result presented here unifies earlier works by Shigekawa [3] and Florit-Nualart [2]. Using the criteria, we investigate existence of smooth densities for the functionals which fail the following properties:

1. Functionals themselves have higher order derivatives,
2. Determinants of the Malliavin covariance matrices have higher order negative moments.

2 Framework and Main Result

Definition 1 (Classical Wiener space). Let $N \in \mathbf{N}$.

$$\begin{aligned} W &:= \{x : [0, 1] \rightarrow \mathbf{R}^N, \text{continuous}, x(0) = 0\}, \\ \|x\|_W &:= \max_{s \in [0, 1]} |x(s)|_{\mathbf{R}^N}, \quad x \in W, \\ \mathcal{H} &:= \{h = (h^1, \dots, h^N) \in W; \\ &\quad h_j \text{ is absolutely continuous w.r.t. Lebesgue measure} \\ &\quad \text{and } \dot{h}_j \text{ is square integrable}\}, \\ (h, k)_{\mathcal{H}} &:= \int_0^1 (\dot{h}(s), \dot{k}(s))_{\mathbf{R}^N} ds, \quad h, k \in \mathcal{H}, \\ P &: \text{Wiener measure.} \end{aligned}$$

Under the measure P , the coordinate process $B_t(x) := x(t)$ defines the N -dimensional Brownian motion starting from the origin.

Definition 2 (\mathcal{H} -derivative). The directional derivative of Wiener functional $F : W \rightarrow \mathbf{R}^d$ along to $h \in \mathcal{H}$ is defined by

$$(DF(x), h)_{\mathcal{H}} := \left. \frac{d}{d\varepsilon} F(x + \varepsilon h) \right|_{\varepsilon=0}.$$

Definition 3 (Sobolev space on Wiener space). Let $n \in \mathbf{N}$ and $1 < p < \infty$.

$$\begin{aligned} \|F\|_{\mathcal{D}^{n,p}(W; \mathbf{R}^d)} &:= \sum_{k=0}^n \|D^k F\|_{\mathcal{L}^p(W; \mathcal{H}^k(\mathcal{H}; \mathbf{R}^d))}, \\ \mathcal{D}^{n,p}(W; \mathbf{R}^d) &:= \overline{\mathcal{P}(W; \mathbf{R}^d)}^{\|\cdot\|_{\mathcal{D}^{n,p}(W; \mathbf{R}^d)}}. \end{aligned}$$

Here, we introduce the notion of generalized locally non-degenerate Wiener functionals and state the main result.

Definition 4 ((A, n, q) -locally non-degenerate). Let $A \subset \mathbf{R}^d$ be an open set, $n \in \mathbf{N}$ and $1 < p, q < \infty$. A functional $F \in \mathcal{D}^{1,p}(W; \mathbf{R}^d)$ is said to be (A, n, q) -locally non-degenerate if there exist $1 < q_1, \dots, q_n < \infty$ and some \mathcal{H}^d -valued functional $U = (U^1, \dots, U^d)$ such that

$$(\text{LND1}) \quad 1/q_1 + \dots + 1/q_n = 1/q,$$

$$(\text{LND2}) \quad U \in \mathcal{D}^{m, q_m}(W; \mathcal{H}^d) \text{ for any } 1 \leq m \leq n,$$

$$(\text{LND3}) \quad 1_{\{F \in A\}}[(DF^j, U^k)_{\mathcal{H}} - \delta_{jk}] = 0 \text{ } \mu\text{-a.s. for any } 1 \leq j, k \leq d.$$

Here, δ_{jk} denotes Kronecker's delta.

Theorem 5. Let $1 < p, q < \infty$ satisfy $1/p + 1/q \leq 1$ and suppose $q > d$. If a functional $F \in \mathcal{D}^{1,p}(W; \mathbf{R}^d)$ is (A, n, q) -locally non-degenerate, then it possesses a density $\rho \in C^{n-1}(A; \mathbf{R})$ on A with respect to Lebesgue measure.

3 Proof of the Main Result

A key lemma to prove the main result is the following integration by parts formula:

Lemma 6. Let $1 < p, q, r < \infty$ satisfy $1/p + 1/q \leq 1$ and suppose that a functional $F \in \mathcal{D}^{1,p}(W; \mathbf{R}^d)$ is (A, n, q) -locally non-degenerate. Then, for any multi-index α with length $1 \leq m \leq n$, there exists $G^\alpha \in \mathcal{D}^{n-m, s_m}(W; \mathbf{R})$ such that

$$\int_W \partial^\alpha \varphi \circ F dP = \int_W (\varphi \circ F) G^\alpha dP \quad \text{for all } \varphi \in C_{\text{cpt}, A}^n(\mathbf{R}^d; \mathbf{R}).$$

Here, $\{s_m\}_{m=1}^n$ is a decreasing sequence defined inductively by

$$\frac{1}{s_m} := \begin{cases} \frac{1}{q_n} + \frac{1}{r}, & m = 1, \\ \frac{1}{q_{n-(m-1)}} + \frac{1}{s_{m-1}}, & 2 \leq m \leq n. \end{cases}$$

From this lemma, we see that there exists $g^\alpha \in \mathcal{L}^q(\mathbf{R}^d, \nu; \mathbf{R})$ such that

$$\int_{\mathbf{R}^d} \partial^\alpha \varphi d\nu = \int_{\mathbf{R}^d} \varphi g^\alpha d\nu \quad \text{for all } \varphi \in C_{\text{cpt}, A}^m(\mathbf{R}^d; \mathbf{R}),$$

where $\nu = P \circ F^{-1}$. This equation and the Sobolev embedding theorem help us to obtain Theorem 5.

4 Application

We study existence of smooth density of δ -dimensional Bessel process r , which is defined by

$$r_t = r_0 + \int_0^t \frac{\delta - 1}{2r_s} ds + B_t,$$

where $r_0 > 0$ and B is one-dimensional Brownian motion. The drift coefficient in this stochastic differential equation has ‘singularity’ at the origin, especially it is not Lipschitz continuous. Thus we have little hope of differentiability of r in general, however, we see:

Proposition 7 ([1]). Let $\delta \geq 2$. We have $r_t \in \mathcal{D}^{1,p}(W; \mathbf{R})$ for any $p > 1$ and

$$Dr_t = \int_0^t 1_{[0, t]}(s) \exp\left(-\int_s^t \frac{\delta - 1}{2r_u^2} du\right) ds.$$

Theorem 8. Let $\delta > 8$ and $1 < q < \delta/6$. Then r_t is $(\mathbf{R}, 1, q)$ -locally non-degenerate and has a continuous density on \mathbf{R} .

To prove non-degeneracy, we choose $U \in \mathcal{D}^{1,q}(W; \mathcal{H})$ as

$$U := \left(1 - \exp\left(-\int_0^t \frac{\delta - 1}{2r_s^2} ds\right)\right)^{-1} \int_0^t 1_{[0, t]}(s) \frac{\delta - 1}{2r_s^2} ds.$$

Using Theorem 5, we obtain the desired result.

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Abstract

Let Γ be a set of functions on the natural numbers. We introduce a new randomness notion called semi Γ -randomness, which is associating with a Γ -indexed test. We prove that weak n -randomness is strictly stronger than semi Δ_n^0 -randomness. Other new results are shown, we investigate the characterization of randomness notions. Moreover, we investigate the relationships between various definitions of randomness. (This is a joint work with Kojiro Higuchi, Takeshi Yamazaki, Kazuyuki Tanaka).

What is random?

Suppose we are given two infinite binary sequences, say

- A sequence of 1's: **111111111111...**
- A sequence obtained by fair coin tosses: **01001110110110...**

Most people will agree that the sequence of coin tosses appear to be more **random** than the sequence of 1's.

- This is due to the fact that the sequence of 1's follows a simple "Law" which we can describe easily.
- The sequence of coin tosses is rather chaotic.

Definition (Martin-Löf, 1966)

- A **Martin-Löf test**, or **ML-test** is a uniformly c.e. sequence $(G_m)_{m \in \mathbb{N}}$ of open sets such that $\forall m \in \mathbb{N} \mu(G_m) \leq 2^{-m}$.
- A sequence $Z \subseteq \mathbb{N}$ **fails** the test if $Z \in \bigcap_m G_m$, otherwise Z **passes** the test.
- Z is **ML-random** if it passes every Martin-Löf test.

Definitions (Kurtz, 1981)

- A **generalized ML-test** is a uniformly c.e. sequence $(G_m)_{m \in \mathbb{N}}$ of open sets such that $\mu(\bigcap_m G_m) = 0$.
- Z is **weakly 2-random** if it passes every generalized ML-test.
- ML-random is also called 1-random. Z is **2-random** if it is 1-random relative to \emptyset' .

Facts

- $2\text{-random} \Rightarrow \text{weakly } 2\text{-random} \Rightarrow \text{ML-random}$.
- The reverse implications fail (Kurtz, Kautz).

Motivation

- This is concerned with the algorithmic notion of randomness such as originally introduced by P. Martin-Löf in 1966.
- One approach is to generalize the Martin-Löf-test by giving the m -th component (a c.e. set of measure at most 2^{-m}) via a function in some function class Γ .
- The main purpose of this post is to give a general framework for such randomness notions.
- To do this, we introduce the notion of semi Γ -randomness.

Definition (Peng, 2011)

A set Z is **Γ -random** if Z is ML-random relative to A for all $A \in \Gamma$. In particular, Γ -randomness is called **\mathbb{L} -randomness** if Γ is the set of low sets.

In fact, it turned out that \mathbb{L} -randomness is equivalent to \emptyset' -Schnorr randomness.

Theorem (Yu, 2012)

For every \emptyset' -Schnorr test $\{U_e^0\}_{e \in \mathbb{N}}$, there is a real z with $z' \leq_T \emptyset'$ such that there is z -Martin-Löf-test $\{V_e^z\}_{e \in \mathbb{N}}$ so that $\bigcap_{e \in \mathbb{N}} V_e^z \supseteq \bigcap_{e \in \mathbb{N}} U_e^0$.

We would like to introduce another characterization of \mathbb{L} -randomness. Then, the next lemma is useful.

Lemma

Let $\Gamma, \Gamma' \subset \mathbb{N}^{\mathbb{N}}$ such that for any $f \in \Gamma$ there is a function $g \in \Gamma \cap \Gamma'$ with $f \leq_{LR} g$. Then, Γ randomness is equivalent to $\Gamma \cap \Gamma'$ randomness.

Proposition

The following are equivalent:

- X is \mathbb{L} randomness,
- X is $\mathbb{L} \cap \mathbb{G}$ randomness,
- X is $\mathbb{L} \cap \mathbf{PA}$ randomness.

Let **MLR** denote the class of ML-random reals.

Theorem

For every \emptyset' -Schnorr test $\{U_e^0\}_{e \in \mathbb{N}}$ and every Π_1^0 class $P \subset \mathbf{MLR}$, there is a low real $z \in P$ such that there is z -Martin-Löf-test $\{V_e^z\}_{e \in \mathbb{N}}$ so that $\bigcap_{e \in \mathbb{N}} V_e^z \supseteq \bigcap_{e \in \mathbb{N}} U_e^0$.

Definition

- Let $\Gamma \subseteq \mathbb{N}^{\mathbb{N}}$. We say that a sequence $\{G_n\}_{n \in \mathbb{N}}$ of c.e. open sets is a **Γ -indexed test** if and only if there exists $f \in \Gamma$ such that $G_n = W_{f(n)}$ for all $n \in \mathbb{N}$ and $\mu(G_n) \leq 2^{-n}$.
- A set $Z \subseteq \mathbb{N}$ **fails** the test if $Z \in \bigcap_n G_n$, otherwise Z **passes** the test.
- A is **semi Γ -random** if it passes every Γ -indexed test.

The following lemma is for separation between weak randomness and semi Γ -randomness.

Lemma

For any $\Gamma, \Gamma' \subset \mathbb{N}^{\mathbb{N}}$ such that Γ randomness is not equivalent to Martin-Löf randomness and Γ' is countable, there exists a semi Γ' random real which is not Γ random.

Main Theorem

If $n > 2$, then weakly n randomness is strictly stronger than semi Δ_n^0 -randomness.

Finally, we consider the case around $n = 2$.

Lemma

There exists a $\Pi_1^0(\emptyset')$ null set P containing a semi \mathbb{L} random element.

Theorem

Weakly 2 randomness is strictly stronger than semi \mathbb{L} -randomness

Future Study

- Study other supper classes of \mathbb{L} -randomness.
- Is the a characterization of \emptyset' -computable randomness via ML randomness?

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Lower bounds of the canonical heights on certain elliptic curves

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Elliptic curves and rational points We denote the set of rational numbers by \mathbb{Q} . An *elliptic curve* E over \mathbb{Q} is a nonsingular curve defined by the equation $y^2 = x^3 + Ax + B$ ($A, B \in \mathbb{Q}$). Then by $E(\mathbb{Q})$, we denote the set of rational points of E with the point at infinity O ($= (\infty, \infty)$). One of the main problems about elliptic curves is to determine $E(\mathbb{Q})$, that is, the rational solution of the above equation. The Mordell–Weil theorem is a basic result about this problem, which states that $E(\mathbb{Q})$ is a finitely generated abelian group. We denote the rank of the abelian group by $\text{rank} E(\mathbb{Q})$.

The canonical height and its computations Heights of a point represent arithmetic complexity of the point. The *canonical height* is one of the heights which is useful to consider $E(\mathbb{Q})$.

Definition 1 For $P = (x, y) \in E(\mathbb{Q})$ we write $x = a/c$ with $\gcd(a, c) = 1$. We define $h(P) = \log \max\{|a|, |c|\}$. We define $\hat{h}(P) = \lim_{n \rightarrow \infty} \frac{1}{4^n} h([2^n]P)$, which is called the *canonical height* of P on $E(\mathbb{Q})$, where $[2^n]P$ means the multiple of P by 2^n in the abelian group $E(\mathbb{Q})$.

A remarkable property of the canonical height is the equality $\hat{h}([m]P) = m^2 \hat{h}(P)$. This means that the canonical heights of points in a basis of $E(\mathbb{Q})$ are relatively small.

The above definition, which is usually used, is brief but not suitable for practical computations. Néron and Tate showed that there is a function λ_p such that $\hat{h}(P) = \sum_{p: \text{prime}, \infty} \lambda_p(P)$, which is called the Néron

local height function. By computing $\lambda_p(P)$ for some p prime number (and ∞), we can estimate the value of the canonical height of P .

We computed a lower bound of the canonical height independent of points.

Theorem 1 Let E/\mathbb{Q} be an elliptic curve defined by $y^2 = x^3 + a_2x^2 + a_4x + a_6$ ($a_2, a_4, a_6 \in \mathbb{Z}$) with the discriminant Δ . Let D be a square-free integer and E_D/\mathbb{Q} the elliptic curve $y^2 = x^3 + a_2Dx^2 + a_4D^2x + a_6D^3$. Assume that Δ is sixth-power-free. Then for $P \in E_D(\mathbb{Q}) \setminus E_D(\mathbb{Q})[2]$,

$$\hat{h}(P) > \frac{1}{4} \log |D| + \frac{1}{16} \log \frac{(1 - |q|)^8}{|q|} + \frac{1}{4} \log \left| \frac{\omega}{2\pi} \right| - \frac{7}{16} \sum_{p|\Delta, p \neq 2} \log p - \frac{5}{12} \log 2,$$

where ω_1 and ω_2 are periods of E such that $\omega_1 > 0$, $\text{Im}(\omega_2) > 0$ and $\text{Re}(\omega_2/\omega_1) = 0$ or $-1/2$, $q = \exp(2\pi i \omega_2/\omega_1)$ and

$$\omega = \begin{cases} \omega_1 & (D > 0) \\ \text{Im}(\omega_2) & (D < 0, \Delta > 0) \\ 2\text{Im}(\omega_2) & (D < 0, \Delta < 0) \end{cases}.$$

Estimates of the canonical height as the above theorem can give an information of the arithmetic of elliptic curves. As an application of it, we can show the following theorem.

Theorem 2 Let $t \in \mathbb{Z}$, $D(t) = t^6 + 4t^4 + 30t^3 + 5t^2 + 54t + 245$, E_D the elliptic curve defined by $y^2 = x^3 + 2D(t)x^2 + 163D(t)^2x + 2205D(t)^3$ and P the point $(D(t)(t^4 + 2t^2 + 12t), D(t)^2(t^3 + t + 3))$ on E_D . We assume that $D(t)$ is square-free. Then P is a primitive if $|t| \geq 2216$. In particular $E_D(\mathbb{Q}) \simeq \langle P \rangle$ if $\text{rank } E_D(\mathbb{Q}) = 1$.